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## Lie Algebras of Finite and Affine Type



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# Lie Algebras of <br> Finite and Affine Type 

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Dedicated to Sandy Green

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## Preface

Lie algebras were originally introduced by S . Lie as algebraic structures used for the study of Lie groups. The tangent space of a Lie group at the identity element has the natural structure of a Lie algebra, called by Lie the infinitesimal group. However, Lie algebras also proved to be of interest in their own right. The finite dimensional simple Lie algebras over the complex field were investigated independently by E. Cartan and W. Killing and the classification of such algebras was achieved during the decade 1890-1900. Basic ideas on the structure and representation theory of these Lie algebras were also contributed at a later stage by H. Weyl. Since then the theory of finite dimensional simple Lie algebras has found many and varied applications both in mathematics and in mathematical physics, to the extent that it is now generally regarded as one of the classical branches of mathematics.

In 1967 V. G. Kac and R. V. Moody independently introduced the Lie algebras now known as Kac-Moody algebras. The finite dimensional simple Lie algebras are examples of Kac-Moody algebras; but the theory of Kac-Moody algebras is much broader, including many infinite dimensional examples. The Kac-Moody theory has developed rapidly since its introduction and has also turned out to have applications in many areas of mathematics, including among others group theory, combinatorics, modular forms, differential equations and invariant theory. It has also proved important in mathematical physics, where it has applications to statistical physics, conformal field theory and string theory. The representation theory of affine Kac-Moody algebras has been particularly useful in such applications.

In view of these applications it seems clear that the theory of Lie algebras, of both finite and affine types, will continue to occupy a central position in mathematics into the twenty-first century. This expectation provides the motivation for the present volume, which aims to give a mathematically rigorous development of those parts of the theory of Lie algebras most relevant
to the understanding of the finite dimensional simple Lie algebras and the Kac-Moody algebras of affine type. A number of books on Lie algebras are confined to the finite dimensional theory, but this seemed too restrictive for the present volume in view of the many current applications of the Kac-Moody theory. On the other hand the Kac-Moody theory needs a prior knowledge of the finite dimensional theory, both to motivate it and to supply many technical details. For this reason I have included an account both of the Cartan-KillingWeyl theory of finite dimensional simple Lie algebras and of the Kac-Moody theory, concentrating particularly on the Kac-Moody algebras of affine type. We work with Lie algebras over the complex field, although any algebraically closed field of characteristic zero would do equally well.

I was introduced to the theory of Lie algebras by an inspiring course of lectures given by Philip Hall at Cambridge University in the late 1950s. I have given a number of lecture courses on finite dimensional Lie algebras at Warwick University, and also two lecture courses on Kac-Moody algebras. The present book has developed as a considerably expanded version of the lecture notes of these courses. The main prerequisite for study of the book is a sound knowledge of linear algebra. I have in fact aimed to make this the sole prerequisite, and to explain from first principles any other techniques which are used in the development.

The most influential book on Kac-Moody algebras is the volume InfiniteDimensional Lie Algebras, third edition (1990), by V. Kac. That formidable treatise contains a development of the Kac-Moody theory presupposing a knowledge of the finite dimensional theory, and includes information on several of the applications. The present volume will not rival Kac' account for experts on Kac-Moody algebras. About half of the theory covered in the 3rd edition of Kac' book has been included. However, for those new to the Kac-Moody theory, our account may be useful in providing a gentler introduction, making use of ideas from the finite dimensional theory developed earlier in the book.

The content of the book can be summarised as follows. The basic definitions of Lie algebras, their subalgebras and ideals, representations and modules, are given in Chapter 1. In Chapter 2 the standard results are proved on the representation theory of soluble and nilpotent Lie algebras. The results on representations of nilpotent Lie algebras are used extensively in the subsequent development. The key idea of a Cartan subalgebra is introduced in Chapter 3, where the existence and conjugacy of Cartan subalgebras are proved. We make use of some ideas from algebraic geometry to prove the conjugacy of Cartan subalgebras. In Chapter 4 the Killing form is introduced and used to describe the Cartan decomposition of a semisimple Lie algebra into root
spaces with respect to a Cartan subalgebra. The well-known example of the special linear Lie algebra is used to illustrate the general ideas. In Chapter 5 the Weyl group is introduced and shown to be a Coxeter group. This leads on to the definition of the Cartan matrix and the Dynkin diagram. The possible Dynkin diagrams and Cartan matrices are classified in Chapter 6, and in Chapter 7 the existence and uniqueness of a semisimple Lie algebra with a given Cartan matrix are proved. In Chapter 8 the finite dimensional simple Lie algebras are discussed individually and their root systems determined.

Chapters 9 to 13 are concerned with the representation theory of finite dimensional semisimple Lie algebras. We begin in Chapter 9 with the introduction of the universal enveloping algebra, of free Lie algebras and of Lie algebras defined by generators and relations. The finite dimensional irreducible modules for semisimple Lie algebras are obtained in Chapter 10 as quotients of infinite dimensional Verma modules with dominant integral highest weight. In Chapter 11 the enveloping algebra is studied in more detail. Its centre is shown to be isomorphic to the algebra of polynomial functions on a Cartan subalgebra invariant under the Weyl group, and to the algebra of polynomial functions on the Lie algebra invariant under the adjoint group. This algebra is shown to be isomorphic to a polynomial algebra. The properties of the Casimir element of the centre of the enveloping algebra are also discussed. These are important in subsequent applications to representation theory. Characters of modules are introduced in Chapter 12, and Weyl's character formula for the irreducible modules is proved. The fundamental irreducible modules for the finite dimensional simple Lie algebras are discussed individually in Chapter 13. Their discussion involves exterior powers of modules, Clifford algebras and spin modules, and contraction maps.

This concludes the development of the structure and representation theory of the finite dimensional Lie algebras. This development has concentrated particularly on the properties necessary to obtain the classification of the simple Lie algebras and their finite dimensional irreducible modules. Among the significant results omitted from our account are Ado's theorem on the existence of a faithful finite dimensional module, the radical splitting theorem of Levi, the theorem of Malcev and Harish-Chandra on the conjugacy of complements to the radical, and the cohomology theory of Lie algebras.

The theory of Kac-Moody algebras is introduced in Chapter 14, where the Kac-Moody algebra associated to a generalised Cartan matrix is defined. In fact there are two slightly different definitions of a Kac-Moody algebra which have been used. There is a definition in terms of generators and relations which appears the more natural, but there is a different definition, given by Kac in his book, which is more convenient when one wishes to show that a
given Lie algebra is a Kac-Moody algebra. I have used the latter definition, but have included a proof that, at least for symmetrisable generalised Cartan matrices, the two definitions are equivalent.

The trichotomy of indecomposable generalised Cartan matrices into those of finite, affine and indefinite types is obtained in Chapter 15. The Kac-Moody algebras of finite type turn out to be precisely the non-trivial finite dimensional simple Lie algebras, and a classification of those of affine type is given. The important special case of symmetrisable Kac-Moody algebras is also introduced. This class includes all those of finite and affine types, and some of those of indefinite type. In Chapter 16 it is shown that symmetrisable algebras have an invariant bilinear form, which plays a key role in the subsequent development. The Weyl group and root system of a Kac-Moody algebra are also discussed. The roots divide into real roots and imaginary roots, and a remarkable theorem of Kac is proved which characterises the set of positive imaginary roots. Kac-Moody algebras of affine type are singled out for more detailed discussion in Chapter 17. In Chapter 18 it is shown how some of them can be realised in terms of a central extension of a loop algebra of a finite dimensional simple Lie algebra, whereas the remainder can be obtained as fixed point subalgebras of these under a twisted graph automorphism.

Chapters 19 and 20 are devoted to the representation theory of Kac-Moody algebras. The representations considered are those from the category $\mathcal{O}$ introduced by Bernstein, Gelfand and Gelfand. In Chapter 19 the irreducible modules in this category are classified, and their characters are obtained in Kac' character formula, a generalisation to the Kac-Moody situation of Weyl's character formula. In Chapter 20 the representations of affine Kac-Moody algebras are discussed. The remarkable identities of I. G. Macdonald are obtained by specialising the denominator of Kac' character formula, interpreted in two different ways; one as an infinite sum and the other as an infinite product. The phenomenon of strings of weights with non-decreasing multiplicities is investigated inside an irreducible module for an affine algebra.

Many of the applications of the representation theory of affine Kac-Moody algebras use the theory of vertex operators. This theory lies beyond the scope of the present volume. However, we have introduced the idea of a vertex operator in Chapter 20 with the aim of encouraging the reader to explore the subject further.

A theory of generalised Kac-Moody algebras was introduced in 1988 by R. Borcherds. These Lie algebras were introduced as part of Borcherds' proof of the Conway-Norton conjectures on the representation theory of the Monster simple group. They are now frequently called Borcherds algebras. In Chapter 21 we have given an account of Borcherds algebras, including the
definition and statements of the main results concerning their structure and representation theory, but detailed proofs are not given. Many of the results on Borcherds algebras are quite similar to those for Kac-Moody algebras, but there are examples of Borcherds algebras which are quite different from Kac-Moody algebras. The best known such example is the Monster Lie algebra, which we describe in the final section.

We conclude with an appendix containing one section for each of the algebras of finite and affine types, in which the most important pieces of information about the algebra concerned are collected.

I would like to express my thanks to Roger Astley of Cambridge University Press for his encouragement to complete the half finished manuscript of this book. This was eventually achieved after I had reached the status of Emeritus Professor, and therefore had more time to devote to it. I would also like to thank my colleague Bruce Westbury for the sustained interest he has shown in this work.

## 1

## Basic concepts

### 1.1 Elementary properties of Lie algebras

A Lie algebra is a vector space $L$ over a field $k$ on which a multiplication

$$
\begin{aligned}
L \times L & \rightarrow L \\
(x, y) & \rightarrow[x y]
\end{aligned}
$$

is defined satisfying the following axioms:
(i) $(x, y) \rightarrow[x y]$ is linear in $x$ and in $y$;
(ii) $[x x]=0 \quad$ for all $x \in L$;
(iii) $[[x y] z]+[[y z] x]+[[z x] y]=0 \quad$ for all $x, y, z \in L$.

Axiom (iii) is called the Jacobi identity.

Proposition 1.1 $[y x]=-[x y]$ for all $x, y \in L$.

Proof. Since $[x+y, x+y]=0$ we have $[x x]+[x y]+[y x]+[y y]=0$. It follows that $[x y]+[y x]=0$, that is $[y x]=-[x y]$.

Proposition 1.1 asserts that multiplication in a Lie algebra is anticommutative.
Now let $H, K$ be subspaces of a Lie algebra $L$. Then $[H K]$ is defined as the subspace spanned by all products $[x y]$ with $x \in H$ and $y \in K$. Each element of [ $H K$ ] is a sum

$$
\left[x_{1} y_{1}\right]+\cdots+\left[x_{r} y_{r}\right]
$$

with $x_{i} \in H, y_{i} \in K$.
Proposition $1.2[H K]=[K H]$ for all subspaces $H, K$ of $L$.

Proof. Let $x \in H, y \in K$. Then $[x y]=[-y, x] \in[K H]$. This shows that $[H K] \subset$ $[K H]$. Similarly we have $[K H] \subset[H K]$ and so we have equality.

Proposition 1.2 asserts that multiplication of subspaces in a Lie algebra is commutative.

Example 1.3 Let $A$ be an associative algebra over $k$. Thus we have a map

$$
\begin{aligned}
A \times A & \rightarrow A \\
(x, y) & \rightarrow x y
\end{aligned}
$$

satisfying the associative law

$$
(x y) z=x(y z) \quad \text { for all } x, y, z \in A
$$

Then $A$ can be made into a Lie algebra by defining the Lie product $[x y]$ by

$$
[x y]=x y-y x
$$

We verify the Lie algebra axioms. Product $[x y]$ is clearly linear in $x$ and in $y$. It is also clear that $[x x]=0$. Finally we check the Jacobi identity. We have

$$
\begin{aligned}
{[[x y] z] } & =(x y-y x) z-z(x y-y x) \\
& =x y z-y x z-z x y+z y x
\end{aligned}
$$

We have similar expressions for $[[y z] x]$ and $[[z x] y]$. Hence

$$
\begin{aligned}
{[[x y] z]+[[y z] x]+[[z x] y]=} & x y z-y x z-z x y+z y x+y z x-z y x-x y z \\
& +x z y+z x y-x z y-y z x+y x z=0
\end{aligned}
$$

The Lie algebra obtained from the associative algebra $A$ in this way will be denoted by $[A]$.

Now let $L$ be a Lie algebra over $k$. A subset $H$ of $L$ is called a subalgebra of $L$ if $H$ is a subspace of $L$ and $[H H] \subset H$. Thus $H$ is itself a Lie algebra under the same operations as $L$.

A subset $I$ of $L$ is called an ideal of $L$ if $I$ is a subspace of $L$ and $[I L] \subset I$. We observe that the latter condition is equivalent to $[L I] \subset I$. Thus there is no distinction between left ideals and right ideals in the theory of Lie algebras. Every ideal is two-sided.

Proposition 1.4 (i) If $H, K$ are subalgebras of $L$ so is $H \cap K$.
(ii) If $H, K$ are ideals of $L$ so is $H \cap K$.
(iii) If $H$ is an ideal of $L$ and $K$ a subalgebra of $L$ then $H+K$ is a subalgebra of $L$.
(iv) If $H, K$ are ideals of $L$ then $H+K$ is an ideal of $L$.

Proof. (i) $H \cap K$ is a subspace of $L$ and $[H \cap K, H \cap K] \subset[H H] \cap[K K] \subset$ $H \cap K$. Thus $H \cap K$ is a subalgebra.
(ii) This time we have $[H \cap K, L] \subset[H L] \cap[K L] \subset H \cap K$. Thus $H \cap K$ is an ideal of $L$.
(iii) $H+K$ is a subspace of $L$. Also $[H+K, H+K] \subset[H H]+[H K]+[K H]+$ $[K K] \subset H+K$, since $[H H] \subset H,[H K] \subset H,[K K] \subset K$. Thus $H+K$ is a subalgebra.
(iv) This time we have $[H+K, L] \subset[H L]+[K L] \subset H+K$. Thus $H+K$ is an ideal of $L$.

We next introduce the idea of a factor algebra. Let $I$ be an ideal of a Lie algebra $L$. Then $I$ is in particular a subspace of $L$ and so we can form the factor space $L / I$ whose elements are the cosets $I+x$ for $x \in L . I+x$ is the subset of $L$ consisting of all elements $y+x$ for $y \in I$.

Proposition 1.5 Let $I$ be an ideal of $L$. Then the factor space $L / I$ can be made into a Lie algebra by defining

$$
[I+x, I+y]=I+[x y] \quad \text { for all } x, y \in L
$$

Proof. We must first show that this definition is unambiguous, that is if $I+x=I+x^{\prime}$ and $I+y=I+y^{\prime}$ then $I+[x y]=I+\left[x^{\prime} y^{\prime}\right]$.

Now $I+x=I+x^{\prime}$ implies that $x=x^{\prime}+i_{1}$ for some $i_{1} \in I$. Similarly $I+y=$ $I+y^{\prime}$ implies $y=y^{\prime}+i_{2}$ for some $i_{2} \in I$. Thus

$$
\begin{aligned}
& I+[x y]=I+\left[x^{\prime}+i_{1}, y^{\prime}+i_{2}\right] \\
& \quad=I+\left[i_{1} y^{\prime}\right]+\left[x^{\prime} i_{2}\right]+\left[i_{1} i_{2}\right]+\left[x^{\prime} y^{\prime}\right] \\
& \quad=I+\left[x^{\prime} y^{\prime}\right]
\end{aligned}
$$

since $\left[i_{1} y^{\prime}\right],\left[x^{\prime} i_{2}\right],\left[i_{1} i_{2}\right]$ all lie in $I$. Thus our multiplication is well defined. We also have

$$
[I+x, I+x]=I+[x x]=I
$$

and the Jacobi identity in $L / I$ clearly follows from the Jacobi identity in $L$.

Now suppose we have two Lie algebras $L_{1}, L_{2}$ over $k$. A map $\theta: L_{1} \rightarrow L_{2}$ is called a homomorphism of Lie algebras if $\theta$ is linear and

$$
\theta[x y]=[\theta x, \theta y] \quad \text { for all } x, y \in L_{1}
$$

The map $\theta: L_{1} \rightarrow L_{2}$ is called an isomorphism of Lie algebras if $\theta$ is a bijective homomorphism of Lie algebras. The Lie algebras $L_{1}, L_{2}$ are said to be isomorphic if there exists an isomorphism $\theta: L_{1} \rightarrow L_{2}$.

Proposition 1.6 Let $\theta: L_{1} \rightarrow L_{2}$ be a homomorphism of Lie algebras. Then the image of $\theta$ is a subalgebra of $L_{2}$, the kernel of $\theta$ is an ideal of $L_{1}$ and $L_{1} / \operatorname{ker} \theta$ is isomorphic to $\operatorname{im} \theta$.

Proof. im $\theta$ is a subspace of $L_{2}$. Moreover for $x, y$ in $L_{1}$ we have

$$
[\theta(x), \theta(y)]=\theta[x y] \in \operatorname{im} \theta
$$

Hence $\operatorname{im} \theta$ is a subalgebra of $L_{2}$.
Now $\operatorname{ker} \theta$ is a subspace of $L_{1}$. Let $x \in \operatorname{ker} \theta$ and $y \in L_{1}$. Then

$$
\theta[x y]=[\theta(x), \theta(y)]=[0, \theta(y)]=0
$$

Hence $[x y] \in \operatorname{ker} \theta$ and so $\operatorname{ker} \theta$ is an ideal of $L_{1}$.
Now let $x, y \in L_{1}$. We consider when $\theta(x)$ is equal to $\theta(y)$. We have

$$
\begin{aligned}
\theta(x)=\theta(y) & \Leftrightarrow \theta(x-y)=0 \Leftrightarrow x-y \in \operatorname{ker} \theta \\
& \Leftrightarrow \operatorname{ker} \theta+x=\operatorname{ker} \theta+y .
\end{aligned}
$$

This shows that there is a bijective map $\theta(x) \rightarrow \operatorname{ker} \theta+x$ between $\operatorname{im} \theta$ and $L_{1} / \operatorname{ker} \theta$. We show this bijection is an isomorphism of Lie algebras. It is clearly linear. Moreover given $x, y, z \in L_{1}$ we have

$$
\begin{aligned}
& {[\theta(x), \theta(y)]=\theta(z) \Leftrightarrow \theta[x y]=\theta(z)} \\
& \quad \Leftrightarrow \operatorname{ker} \theta+[x y]=\operatorname{ker} \theta+z \\
& \quad \Leftrightarrow[\operatorname{ker} \theta+x, \operatorname{ker} \theta+y]=\operatorname{ker} \theta+z
\end{aligned}
$$

Thus the bijection preserves Lie multiplication, so is an isomorphism of Lie algebras.

Proposition 1.7 Let I be an ideal of $L$ and $H$ a subalgebra of $L$. Then
(i) $I$ is an ideal of $I+H$.
(ii) $I \cap H$ is an ideal of $H$.
(iii) $(I+H) / I$ is isomorphic to $H /(I \cap H)$.

Proof. We recall from Proposition 1.4 that $I \cap H$ and $I+H$ are subalgebras. We have $[I, I+H] \subset[I L] \subset I$, thus $I$ is an ideal of $I+H$. Also $[I \cap H, H] \subset$ $[I H] \cap[H H] \subset I \cap H$, thus $I \cap H$ is an ideal of $H$.

Let $\theta: H \rightarrow(I+H) / I$ be defined by $\theta(x)=I+x$. This is clearly a linear map, and is also evidently a homomorphism of Lie algebras. It is surjective since each element of $(I+H) / I$ has form $I+x$ for some $x \in H$. Finally its kernel is the set of $x \in H$ for which $I+x=I$, that is $I \cap H$. Thus $(I+H) / I$ is isomorphic to $H /(I \cap H)$ by Proposition 1.6.

### 1.2 Representations and modules

Let $M_{n}(k)$ be the associative algebra of all $n \times n$ matrices over the field $k$ and let $\left[M_{n}(k)\right]$ be the corresponding Lie algebra. This is often called the general linear Lie algebra of degree $n$ over $k$ and we write

$$
\mathfrak{g l}_{n}(k)=\left[M_{n}(k)\right] .
$$

We have $\operatorname{dimgl} l_{n}(k)=n^{2}$.
A representation of a Lie algebra $L$ over $k$ is a homomorphism of Lie algebras

$$
\rho: L \rightarrow \mathfrak{g l}_{n}(k)
$$

for some $n$, and $\rho$ is called a representation of degree $n$. Two representations $\rho, \rho^{\prime}$ of degree $n$ are called equivalent if there exists a non-singular $n \times n$ matrix $T$ such that

$$
\rho^{\prime}(x)=T^{-1} \rho(x) T \quad \text { forall } x \in L
$$

A left L-module is a vector space $V$ over $k$ together with a multiplication

$$
\begin{aligned}
L \times V & \rightarrow V \\
(x, v) & \rightarrow x v
\end{aligned}
$$

satisfying the axioms:
(i) $(x, v) \rightarrow x v$ is linear in $x$ and in $v$;
(ii) $[x y] v=x(y v)-y(x v) \quad$ for all $x, y \in L$ and $v \in V$.

Suppose $V$ is a finite dimensional $L$-module. Let $e_{1}, \ldots, e_{n}$ be a basis of $V$. Let

$$
x e_{j}=\sum_{i} \rho_{i j}(x) e_{i}
$$

with $\rho_{i j}(x) \in k$ and let $\rho(x)=\left(\rho_{i j}(x)\right)$. Then $\rho$ is a representation of $L$. For we have

$$
\begin{aligned}
{[x y] e_{j} } & =x\left(y e_{j}\right)-y\left(x e_{j}\right) \\
& =x\left(\sum_{k} \rho_{k j}(y) e_{k}\right)-y\left(\sum_{k} \rho_{k j}(x) e_{k}\right) \\
& =\sum_{k} \rho_{k j}(y) x e_{k}-\sum_{k} \rho_{k j}(x) y e_{k} \\
& =\sum_{k} \rho_{k j}(y)\left(\sum_{i} \rho_{i k}(x) e_{i}\right)-\sum_{k} \rho_{k j}(x)\left(\sum_{i} \rho_{i k}(y) e_{i}\right) \\
& =\sum_{i}\left(\sum_{k}\left(\rho_{i k}(x) \rho_{k j}(y)-\rho_{i k}(y) \rho_{k j}(x)\right)\right) e_{i} \\
& =\sum_{i}(\rho(x) \rho(y)-\rho(y) \rho(x))_{i j} e_{i}
\end{aligned}
$$

Thus $\rho[x y]=\rho(x) \rho(y)-\rho(y) \rho(x)=[\rho(x), \rho(y)]$ and $\rho$ is a representation of $L$.

Suppose now we take a second basis $f_{1}, \ldots, f_{n}$ of $V$. Let $\rho^{\prime}$ be the representation of $L$ obtained from this basis. Then $\rho^{\prime}$ is equivalent to $\rho$. For there exists a non-singular $n \times n$ matrix $T$ such that

$$
f_{j}=\sum_{i} T_{i j} e_{i}
$$

Thus we have

$$
x f_{j}=\sum_{k} T_{k j} x e_{k}=\sum_{k} T_{k j}\left(\sum_{i} \rho_{i k}(x) e_{i}\right)=\sum_{i}\left(\sum_{k} \rho_{i k}(x) T_{k j}\right) e_{i}
$$

On the other hand

$$
x f_{j}=\sum_{k} \rho_{k j}^{\prime}(x) f_{k}=\sum_{k} \rho_{k j}^{\prime}(x)\left(\sum_{i} T_{i k} e_{i}\right)=\sum_{i}\left(\sum_{k} T_{i k} \rho_{k j}^{\prime}(x)\right) e_{i}
$$

It follows that $\rho(x) T=T \rho^{\prime}(x)$, that is $\rho^{\prime}(x)=T^{-1} \rho(x) T$ for all $x \in L$. Hence the representation $\rho^{\prime}$ is equivalent to $\rho$.

Example 1.8 $L$ is itself a left $L$-module.
The left action of $L$ on $L$ is defined as $x \cdot y=[x y]$. Then we have

$$
[[x y] z]=[x[y z]]-[y[x z]]
$$

which is a consequence of the Jacobi identity. This shows that $L$ is a left $L$-module. This is called the adjoint module. We define ad $x: L \rightarrow L$ by

$$
\text { ad } x \cdot y=[x y] \quad \text { for } x, y \in L
$$

Then we have

$$
\operatorname{ad}[x y]=\operatorname{ad} x \text { ad } y-\operatorname{ad} y \text { ad } x .
$$

Now let $V$ be a left $L$-module, $U$ be a subspace of $V$ and $H$ a subspace of $L$. We define $H U$ to be the subspace of $V$ spanned by all elements of the form $x u$ for $x \in H, u \in U$.

A submodule of $V$ is a subspace $U$ of $V$ such that $L U \subset U$. In particular $V$ is a submodule of $V$ and the zero subspace $O=\{0\}$ is a submodule of $V$. A proper submodule of $V$ is a submodule distinct from $V$ and $O$.

An $L$-module $V$ is called irreducible if it has no proper submodules. $V$ is called completely reducible if it is a direct sum of irreducible submodules. $V$ is called indecomposable if $V$ cannot be written as a direct sum of two proper submodules. Of course every irreducible $L$-module is indecomposable, but the converse need not be true.

We may also define right $L$-modules, but we shall mainly work with left $L$-modules, and $L$-modules will be assumed to be left $L$-modules unless otherwise stated.

### 1.3 Abelian, nilpotent and soluble Lie algebras

A Lie algebra $L$ is abelian if $[L L]=O$. Thus $[x y]=0$ for all $x, y \in L$ when $L$ is abelian.

Given any Lie algebra $L$ we define the powers of $L$ by

$$
L^{1}=L, \quad L^{n+1}=\left[L^{n} L\right] \quad \text { for } n \geq 1
$$

Thus $L$ is abelian if and only if $L^{2}=O$.

Proposition 1.9 $L^{n}$ is an ideal of L. Also

$$
L=L^{1} \supset L^{2} \supset L^{3} \supset \cdots .
$$

Proof. We first observe that if $I, J$ are ideals of $L$ then $[I J]$ is also an ideal of $L$. For let $x \in I, y \in J, z \in L$. Then

$$
[[x y] z]=[x[y z]]-[y[x z]] \in[I J] .
$$

It follows that $L^{n}$ is an ideal of $L$ for each $n>0$. Thus we have

$$
L^{n+1}=\left[L^{n} L\right] \subset L^{n}
$$

A Lie algebra $L$ is called nilpotent if $L^{n}=O$ for some $n \geq 1$. Thus every abelian Lie algebra is nilpotent. It is clear that every subalgebra and every factor algebra of a nilpotent Lie algebra are nilpotent.

We now consider a different kind of powers of $L$. We define

$$
L^{(0)}=L, \quad L^{(n+1)}=\left[L^{(n)}, L^{(n)}\right] \quad \text { for } n \geq 0
$$

Proposition 1.10 $L^{(n)}$ is an ideal of L. Also

$$
L=L^{(0)} \supset L^{(1)} \supset L^{(2)} \supset \cdots
$$

Proof. $L^{(n)}$ is an ideal of $L$ since the product of two ideals is an ideal. Also

$$
L^{(n+1)}=\left[L^{(n)}, L^{(n)}\right] \subset L^{(n)}
$$

A Lie algebra $L$ is called soluble if $L^{(n)}=O$ for some $n \geq 0$.
Proposition 1.11 (a) $\left[L^{m} L^{n}\right] \subset L^{m+n}$ for all $m, n \geq 1$. (b) $L^{(n)} \subset L^{2^{n}}$ for all $n \geq 0$. (c) Every nilpotent Lie algebra is soluble.

Proof. (a). We use induction on $n$. The result is clear if $n=1$. Suppose it is true for $n=r$. Then

$$
\begin{aligned}
{\left[L^{m} L^{r+1}\right] } & =\left[L^{m}\left[L^{r} L\right]\right]=\left[\left[L^{r} L\right] L^{m}\right] \\
& \subset\left[\left[L L^{m}\right] L^{r}\right]+\left[\left[L^{m} L^{r}\right] L\right] \quad \text { by the Jacobi identity } \\
& \subset\left[L^{m+1} L^{r}\right]+\left[\left[L^{m} L^{r}\right] L\right] \\
& \subset L^{m+r+1} \quad \text { by inductive hypothesis. }
\end{aligned}
$$

Thus the result holds for $n=r+1$, so for all $n$.
(b). We again use induction on $n$. The result is clear if $n=1$. Suppose it is true for $n=r$. Then

$$
L^{(r+1)}=\left[L^{(r)} L^{(r)}\right] \subset\left[L^{2^{r}} L^{2^{r}}\right] \subset L^{2^{r+1}}
$$

by (a). Thus the result holds for $n=r+1$, so for all $n$.
(c). Suppose $L$ is nilpotent. Then $L^{2^{n}}=O$ for $n$ sufficiently large. Hence $L^{(n)}=O$ by (b) and so $L$ is soluble.

It is clear that every subalgebra and every factor algebra of a soluble Lie algebra are soluble.

Proposition 1.12 Suppose I is an ideal of $L$ and both $I$ and $L / I$ are soluble. Then $L$ is soluble.

Proof. Since $L / I$ is soluble we have $(L / I)^{(n)}=O$ for some $n$. This implies $L^{(n)} \subset I$. Since $I$ is soluble we have $I^{(m)}=O$ for some $m$. Hence

$$
L^{(n+m)}=\left(L^{(n)}\right)^{(m)} \subset I^{(m)}=O
$$

and so $L$ is soluble.

Proposition 1.13 Every finite dimensional Lie algebra L contains a unique maximal soluble ideal $R$. Also $L / R$ contains no non-zero soluble ideal.

Proof. Let $I, J$ be soluble ideals of $L$. Then $I+J$ is also an ideal of $L$ and $(I+J) / I$ is isomorphic to $J /(I \cap J)$ by Proposition 1.7. Now $J$ is soluble, thus $J /(I \cap J)$ is soluble and so $(I+J) / I$ is soluble. Since $I$ is soluble we see that $I+J$ is soluble by Proposition 1.12. Thus the sum of two soluble ideals of $L$ is a soluble ideal. It follows that $L$ has a unique maximal soluble ideal $R$.

If $I / R$ is a soluble ideal of $L / R$ then $I$ is a soluble ideal of $L$ by Proposition 1.12. Hence $I=R$ and $I / R=O$.

The ideal $R$ is called the soluble radical of $L$. A Lie algebra $L$ is called semisimple if $R=O$. Thus $L$ is semisimple if and only if $L$ has no non-zero soluble ideal.
$L$ is called simple if $L$ has no proper ideal, that is no ideal other than $L$ and $O$.

Suppose $L$ is a Lie algebra of dimension 1 over $k$. Then $L$ has a basis $\{x\}$ with 1 element. Since $[x x]=0$ we have $L^{2}=O$. Thus $L$ is abelian. We see that any two 1-dimensional Lie algebras over $k$ are isomorphic. Of course any such Lie algebra is simple, because $L$ has no proper subspaces. The 1-dimensional Lie algebra is called the trivial simple Lie algebra. A non-trivial simple Lie algebra is a simple Lie algebra $L$ with $\operatorname{dim} L>1$.

Proposition 1.14 Each non-trivial simple Lie algebra is semisimple.

Proof. Suppose $L$ is simple but not semisimple. Then the soluble radical $R$ satisfies $R \neq O$. Since $R$ is an ideal of $L$ this implies $R=L$. Thus $L$ is soluble.

Hence $L^{(n)}=O$ for some $n \geq 0$. This implies that $L^{(1)} \neq L$ since $L^{(1)}=L$ would imply $L^{(n)}=L$ for all $n$. Now $L^{(1)}$ is an ideal of $L$, hence $L^{(1)}=O$ since $L$ is simple. Thus $[L L]=O$. But then every subspace of $L$ is an ideal of $L$. Since $L$ is simple $L$ has no proper subspaces, so $\operatorname{dim} L=1$. Thus the only simple Lie algebra which is not semisimple is the trivial simple Lie algebra.

## 2

# Representations of soluble and nilpotent Lie algebras 

### 2.1 Representations of soluble Lie algebras

We shall now and subsequently take the base field $k$ to be the field $\mathbb{C}$ of complex numbers. We shall also assume until further notice that $L$ is a finite dimensional Lie algebra over $\mathbb{C}$, although at a later stage we shall also consider infinite dimensional Lie algebras.

We first consider 1-dimensional representations of a Lie algebra $L$. A 1-dimensional representation is a linear map $\rho: L \rightarrow \mathbb{C}$ such that $\rho[x y]=$ $[\rho(x), \rho(y)]$ for all $x, y \in L$.

Lemma 2.1 A linear map $\rho: L \rightarrow \mathbb{C}$ is a 1-dimensional representation of $L$ if and only if $\rho$ vanishes on $L^{2}$.

Proof. Suppose $\rho$ is a representation. Then for $x, y \in L$ we have

$$
\rho[x y]=[\rho(x), \rho(y)]=\rho(x) \rho(y)-\rho(y) \rho(x)=0 .
$$

Hence $\rho$ vanishes on $L^{2}$.
Conversely suppose that $\rho$ vanishes on $L^{2}$. Then

$$
\rho[x y]=0=[\rho(x), \rho(y)]
$$

and so $\rho$ is a representation of $L$.
We shall now prove a theorem of Lie which shows that any irreducible representation of a soluble Lie algebra is 1-dimensional.

Theorem 2.2 (Lie's theorem). Let L be a soluble Lie algebra and $V$ be a finite dimensional irreducible L-module. Then $\operatorname{dim} V=1$.

Proof. Since $L$ is soluble we have $L^{2} \neq L$. Let $I$ be a subspace of $L$ such that $I \supset L^{2}$ and $\operatorname{dim} I=\operatorname{dim} L-1$. Then $I$ is an ideal of $L$ since

$$
[I L] \subset[L L]=L^{2} \subset I
$$

Thus $I$ is an ideal of $L$ of codimension 1.
We shall prove Lie's theorem by induction on $\operatorname{dim} L$. Suppose $\operatorname{dim} L=1$ and $V$ be an irreducible $L$-module. Let $L=\mathbb{C} x$ and $v$ be an eigenvector of $x$ in $V$. Then $\mathbb{C} v$ is an $L$-submodule of $V$. Since $V$ is irreducible we have $V=\mathbb{C} v$ and $\operatorname{dim} V=1$.

Now suppose $\operatorname{dim} L>1$ and $V$ is an irreducible $L$-module. We may regard $V$ as an $I$-module. Then $V$ contains an irreducible $I$-submodule $W$ and we may assume $\operatorname{dim} W=1$ by induction. Let $w$ be a non-zero vector in $W$. Then

$$
y w=\lambda(y) w \quad \text { for all } y \in I
$$

where $\lambda$ is the 1 -dimensional representation of $I$ given by $W$. Let

$$
U=\{u \in V ; y u=\lambda(y) u \quad \text { for all } y \in I\}
$$

Then we have

$$
O \neq W \subset U \subset V
$$

We shall show that $U$ is an $L$-submodule of $V$. Let $u \in U, x \in L$. Then

$$
y(x u)=x(y u)-[x y] u=\lambda(y) x u-\lambda([x y]) u
$$

since $[x y] \in I$. We shall show $\lambda([x y])=0$. Once we know this we have $x u \in U$ and so $U$ is an $L$-submodule. Since $V$ is irreducible we have $U=V$. Hence

$$
y v=\lambda(y) v \quad \text { for all } v \in V, y \in I
$$

Since $\operatorname{dim} I=\operatorname{dim} L-1$ we can write $L=I \oplus \mathbb{C} x$, a direct sum of subspaces. Let $v$ be an eigenvector for $x$ on $V$. Then $\mathbb{C} v$ is an $L$-submodule of $V$, being invariant under the action of both $I$ and $x$. Since $V$ is irreducible we have $V=\mathbb{C} v$ and so $\operatorname{dim} V=1$.

In order to complete the proof we must show that $\lambda([x y])=0$ for all $x \in L$, $y \in I$. In fact it is sufficient to prove this for the element $x$ chosen above such that $L=I \oplus \mathbb{C} x$.

Let $u$ be any non-zero element of $U$. We write

$$
v_{0}=u, \quad v_{1}=x u, \quad v_{2}=x(x u), \ldots
$$

We have $v_{0}, v_{1}, v_{2}, \ldots \in V$ and so there exists $p \geq 0$ such that $v_{0}, v_{1}, \ldots, v_{p}$ are linearly independent and $v_{p+1}$ is a linear combination of these. Consider the subspace $\left\langle v_{0}, v_{1}, \ldots, v_{p}\right\rangle$ of $V$ spanned by these vectors. This subspace
is invariant under the action of $x$. We consider the effect on this subspace of elements $y \in I$. We have

$$
y v_{0}=y u=\lambda(y) u=\lambda(y) v_{0} .
$$

We shall show

$$
y v_{i}=\lambda(y) v_{i}+\text { a linear combination of } v_{0}, \ldots, v_{i-1} .
$$

This is true for $i=0$. Assuming it for $v_{i-1}$ we have

$$
\begin{aligned}
y v_{i}= & y\left(x v_{i-1}\right)=x\left(y v_{i-1}\right)-[x y] v_{i-1} \\
= & x\left(\lambda(y) v_{i-1}+\text { a linear combination of } v_{0}, \ldots, v_{i-2}\right) \\
& -\left(\text { a linear combination of } v_{0}, \ldots, v_{i-1}\right) \\
= & \lambda(y) v_{i}+\text { a linear combination of } v_{0}, \ldots, v_{i-1} .
\end{aligned}
$$

Thus the subspace $\left\langle v_{0}, v_{1}, \ldots, v_{p}\right\rangle$ is invariant under the action of $y$ for all $y \in I$, as well as being invariant under $x$. Hence it is an $L$-submodule of $V$. Since $V$ is irreducible we have

$$
V=\left\langle v_{0}, v_{1}, \ldots, v_{p}\right\rangle
$$

Now $[x y] \in I$ and we see from the above description of the action of $I$ that

$$
\operatorname{trace}_{V}[x y]=(p+1) \lambda([x y])
$$

Thus we have $(p+1) \lambda([x y])=\operatorname{trace}_{V}[x y]=\operatorname{trace}_{V} x y-\operatorname{trace}_{V} y x=0$, since $\operatorname{trace}_{V} x y=\operatorname{trace}_{V} y x$. Hence $\lambda([x y])=0$ and the proof is complete.

Corollary 2.3 Let $L$ be soluble and $V$ be a finite dimensional L-module. Then a basis can be chosen for $V$ with respect to which we obtain a matrix representation $\rho$ of $L$ of the form

$$
\rho(x)=\left(\begin{array}{cccccc}
* & & & & & \\
0 & * & & & * \\
0 & & \cdot & & & \\
\cdot & & 0 & \cdot & & \\
\cdot & & & & * & \\
0 & \cdot & \cdot & 0 & 0 & *
\end{array}\right) \quad \text { for all } x \in L
$$

Thus the matrices representing elements of $L$ are all of triangular form.

Corollary 2.4 Let L be a soluble Lie algebra with $\operatorname{dim} L=n$. Then $L$ has a chain of ideals

$$
O=I_{0} \subset I_{1} \subset \cdots \subset I_{n-1} \subset I_{n}=L
$$

with $\operatorname{dim} I_{r}=r$.

Proof. We apply Theorem 2.2 to the adjoint $L$-module $L$. The submodules of $L$ are the ideals of $L$. By taking a maximal chain of submodules we obtain ideals of $L$ with the required property.

### 2.2 Representations of nilpotent Lie algebras

When $L$ is a nilpotent Lie algebra we can obtain even stronger results about its representations. Moreover these results on representations of nilpotent Lie algebras play a crucial role in the understanding of semisimple Lie algebras, which we shall deal with subsequently. We begin by recalling results from linear algebra related to the Jordan canonical form. Any $n \times n$ matrix over $\mathbb{C}$ is similar to a diagonal sum of Jordan block matrixes of form

$$
\left(\begin{array}{ccccccc}
\lambda & 1 & & & & \\
& \lambda & 1 & & & 0 & \\
& & \cdot & \cdot & & & \\
& & & \cdot & \cdot & & \\
& & & & \cdot & 1 & \\
& 0 & & & & \lambda & 1 \\
& & & & & & \lambda
\end{array}\right)
$$

In a similar way any linear transformation $\theta: V \rightarrow V$ on a finite dimensional vector space $V$ over $\mathbb{C}$ gives rise to a decomposition of $V$ as in the following proposition.

Proposition 2.5 Let $\theta: V \rightarrow V$ be a linear map with characteristic polynomial

$$
\chi(t)=\left(t-\lambda_{1}\right)^{m_{1}}\left(t-\lambda_{2}\right)^{m_{2}} \ldots\left(t-\lambda_{r}\right)^{m_{r}}
$$

where $\lambda_{1}, \ldots, \lambda_{r}$ are the distinct eigenvalues of $\theta$ and $m_{1}, \ldots, m_{r}$ are their multiplicities. Let $V_{i}$ be the set of all $v \in V$ annihilated by some power of $\theta-\lambda_{i} 1$. Then we have

$$
V=V_{1} \oplus V_{2} \oplus \cdots \oplus V_{r}
$$

Moreover $\operatorname{dim} V_{i}=m_{i}, \theta\left(V_{i}\right) \subset V_{i}$ and the characteristic polynomial of $\theta$ on $V_{i}$ is $\left(t-\lambda_{i}\right)^{m_{i}}$.

Proof. Although this is a standard result from linear algebra we shall prove it in view of its importance for the theory of Lie algebras.

We begin by showing that $V_{i}$ is equal to $W_{i}=\left\{v \in V ;\left(\theta-\lambda_{i} 1\right)^{m_{i}} v=0\right\}$. It is clear that $W_{i} \subset V_{i}$. So let $v \in V_{i}$. Then

$$
\left(\theta-\lambda_{i} 1\right)^{N} v=0 \quad \text { for some } N
$$

We may choose $N \geq m_{i}$. Also

$$
\prod_{j=1}^{r}\left(\theta-\lambda_{j} 1\right)^{m_{j}} v=0
$$

for, by the Cayley-Hamilton theorem, $\theta$ satisfies its own characteristic equation. Now the polynomials

$$
\left(t-\lambda_{i}\right)^{N}, \quad \prod_{j=1}^{r}\left(t-\lambda_{j}\right)^{m_{j}}
$$

have highest common factor $\left(t-\lambda_{i}\right)^{m_{i}}$. Thus there exist polynomials $p(t), q(t) \in \mathbb{C}[t]$ such that

$$
\left(t-\lambda_{i}\right)^{m_{i}}=p(t)\left(t-\lambda_{i}\right)^{N}+q(t) \prod_{j=1}^{r}\left(t-\lambda_{j}\right)^{m_{j}}
$$

Hence

$$
\left(\theta-\lambda_{i} 1\right)^{m_{i}} v=p(\theta)\left(\theta-\lambda_{i} 1\right)^{N} v+q(\theta) \prod_{j=1}^{r}\left(\theta-\lambda_{j} 1\right)^{m_{j}} v=0
$$

Thus $v \in W_{i} \quad$ and $\quad V_{i}=W_{i}$.
We next show that $V=V_{1} \oplus \cdots \oplus V_{r}$. Let $f_{i}(t)=\left(t-\lambda_{1}\right)^{m_{1}} \ldots\left(t-\lambda_{i-1}\right)^{m_{i-1}}$ $\left(t-\lambda_{i+1}\right)^{m_{i+1}} \ldots\left(t-\lambda_{r}\right)^{m_{r}}$. Then the polynomials $f_{1}(t), \ldots, f_{r}(t)$ have highest common factor 1 . Thus there exist polynomials $p_{1}(t), \ldots, p_{r}(t) \in \mathbb{C}[t]$ with $\sum_{i} f_{i}(t) p_{i}(t)=1$.

Let $v \in V$. Then $v=\sum_{i} f_{i}(\theta) p_{i}(\theta) v$. Let $v_{i}=f_{i}(\theta) p_{i}(\theta) v$. Then

$$
\left(\theta-\lambda_{i} 1\right)^{m_{i}} v_{i}=\chi(\theta)\left(p_{i}(\theta) v\right)=0
$$

by the Cayley-Hamilton theorem. Thus $v_{i} \in V_{i}$. Hence $v=v_{1}+\cdots+v_{r}$ with $v_{i} \in V_{i}$, and so $V=V_{1}+\cdots+V_{r}$.

In order to show the sum is direct we must prove

$$
V_{i} \cap\left(V_{1}+\cdots+V_{i-1}+V_{i+1}+\cdots+V_{r}\right)=O
$$

Now the polynomials $\left(t-\lambda_{i}\right)^{m_{i}}$ and $f_{i}(t)$ have highest common factor 1 , thus there exist $p(t), q(t) \in \mathbb{C}[t]$ with

$$
p(t)\left(t-\lambda_{i}\right)^{m_{i}}+q(t) f_{i}(t)=1
$$

Let $v \in V_{i} \cap\left(V_{1}+\cdots+V_{i-1}+V_{i+1}+\cdots+V_{r}\right)$. Since $v \in V_{i}$ we have

$$
\left(\theta-\lambda_{i} 1\right)^{m_{i}} v=0
$$

Since $v \in V_{1}+\cdots+V_{i-1}+V_{i+1}+\cdots+V_{r}$ we have

$$
f_{i}(\theta) v=0
$$

Hence $v=p(\theta)\left(\theta-\lambda_{i} 1\right)^{m_{i}} v+q(\theta) f_{i}(\theta) v=0$. Thus we have shown that

$$
V=V_{1} \oplus \cdots \oplus V_{r}
$$

We next observe that $\theta$ acts on each $V_{i}$. For let $v \in V_{i}$. Then

$$
\left(\theta-\lambda_{i} 1\right)^{m_{i}} \theta v=\theta\left(\theta-\lambda_{i} 1\right)^{m_{i}} v=\theta(0)=0
$$

thus $\theta v \in V_{i}$.
We next show that the only eigenvalue of $\theta: V_{i} \rightarrow V_{i}$ is $\lambda_{i}$. Suppose if possible that $\lambda_{j}$ is an eigenvalue for some $j \neq i$ and let $v \in V_{i}$ be an eigenvector for $\lambda_{j}$. Then $v \neq 0,\left(\theta-\lambda_{i} 1\right)^{m_{i}} v=0$ and $\left(\theta-\lambda_{j} 1\right) v=0$. But the polynomials $\left(t-\lambda_{i}\right)^{m_{i}}$ and $t-\lambda_{j}$ have highest common factor 1 so there exist $p(t), q(t) \in$ $\mathbb{C}[t]$ with

$$
p(t)\left(t-\lambda_{i}\right)^{m_{i}}+q(t)\left(t-\lambda_{j}\right)=1
$$

Hence $v=p(\theta)\left(\theta-\lambda_{i} 1\right)^{m_{i}} v+q(\theta)\left(\theta-\lambda_{j} 1\right) v=0$, a contradiction. So all eigenvalues of $\theta_{i}: V_{i} \rightarrow V_{i}$ are equal to $\lambda_{i}$. It follows that $\operatorname{dim} V_{i} \leq m_{i}$ since $m_{i}$ is the multiplicity of eigenvalue $\lambda_{i}$ on V . But

$$
\operatorname{dim} V=\operatorname{dim} V_{1}+\cdots+\operatorname{dim} V_{r}=m_{1}+\cdots+m_{r}
$$

It follows that $\operatorname{dim} V_{i}=m_{i}$. Finally the characteristic polynomial of $\theta$ on $V_{i}$ is $\left(t-\lambda_{i}\right)^{m_{i}}$.

The subspace $V_{i}$ is called the generalised eigenspace of $V$ with eigenvalue $\lambda_{i}$. Thus the ordinary eigenspace of $\lambda_{i}$ lies in the generalised eigenspace. It is not in general true that $V$ is the direct sum of its eigenspaces with respect to its different eigenvalues, but Proposition 2.5 shows that this result is true if the eigenspaces are replaced by the generalised eigenspaces.

The relevance of the decomposition into generalised eigenspaces for the representations of nilpotent Lie algebras is shown by the following theorem.

Theorem 2.6 Let L be a nilpotent Lie algebra and $V$ be an L-module. Let $y \in L$ and $\rho(y): V \rightarrow V$ be the map $v \rightarrow y v$. Then the generalised eigenspaces $V_{i}$ of $V$ associated with $\rho(y)$ are all submodules of $V$.

Before proving this theorem we need a preliminary result.

Proposition 2.7 Let $L$ be a Lie algebra and $V$ be an L-module. Let $v \in V$, $x, y \in L$ and $\alpha, \beta \in \mathbb{C}$. Then

$$
(\rho(y)-(\alpha+\beta) 1)^{n} x v=\sum_{i=0}^{n}\binom{n}{i}\left((\operatorname{ad} y-\beta 1)^{i} x\right)\left((\rho(y)-\alpha 1)^{n-i} v\right)
$$

Proof. We use induction on $n$. The result is clear when $n=0$. We assume it for $n=r$. We write

$$
x_{i}=(\operatorname{ad} y-\beta 1)^{i} x \in L
$$

Then we have

$$
(\rho(y)-(\alpha+\beta) 1)^{r+1} x v=(\rho(y)-(\alpha+\beta) 1) \sum_{i=0}^{r}\binom{r}{i} \rho\left(x_{i}\right)(\rho(y)-\alpha 1)^{r-i} v
$$

Now

$$
\begin{aligned}
(\rho(y)-(\alpha+\beta) 1) \rho\left(x_{i}\right) & =\rho\left(\left[y x_{i}\right]\right)+\rho\left(x_{i}\right) \rho(y)-(\alpha+\beta) \rho\left(x_{i}\right) \\
& =\rho\left((\operatorname{ad} y-\beta 1) x_{i}\right)+\rho\left(x_{i}\right)(\rho(y)-\alpha 1) \\
& =\rho\left(x_{i+1}\right)+\rho\left(x_{i}\right)(\rho(y)-\alpha 1)
\end{aligned}
$$

Hence

$$
\begin{aligned}
& (\rho(y)-(\alpha+\beta) 1)^{r+1} x v \\
& =\sum_{i=0}^{r}\binom{r}{i} \rho\left(x_{i+1}\right)(\rho(y)-\alpha 1)^{r-i} v+\sum_{i=0}^{r}\binom{r}{i} \rho\left(x_{i}\right)(\rho(y)-\alpha 1)^{r+1-i} v \\
& = \\
& =\sum_{i=0}^{r+1}\binom{r}{i-1} \rho\left(x_{i}\right)(\rho(y)-\alpha 1)^{r+1-i} v+\sum_{i=0}^{r+1}\binom{r}{i} \rho\left(x_{i}\right)(\rho(y)-\alpha 1)^{r+1-i} v \\
& \quad \quad\left(\text { interpreting }\binom{r}{-1}=0 \text { and }\binom{r}{r+1}=0\right) \\
& = \\
& \sum_{i=0}^{r+1}\binom{r+1}{i}\left((\operatorname{ad} y-\beta 1)^{i} x\right)\left((\rho(y)-\alpha 1)^{r+1-i} v\right) .
\end{aligned}
$$

This completes the induction.

Proof of Theorem 2.6. Let $v \in V_{i}, x, y \in L$. Then

$$
\left(\rho(y)-\lambda_{i} 1\right)^{n} x v=\sum_{j=0}^{n}\binom{n}{j}\left((\operatorname{ad} y)^{j} x\right)\left(\left(\rho(y)-\lambda_{i} 1\right)^{n-j} v\right)
$$

by Proposition 2.7 with $\alpha=\lambda_{i}, \beta=0$. Since $v \in V_{i},\left(\rho(y)-\lambda_{i} 1\right)^{n-j} v=0$ if $n-j$ is sufficiently large. Since $L$ is nilpotent $(\operatorname{ad} y)^{j} x=0$ if $j$ is sufficiently large. Thus $\left(\rho(y)-\lambda_{i} 1\right)^{n} x v=0$ if $n$ is sufficiently large. Hence $x v \in V_{i}$ and so $V_{i}$ is a submodule of $V$.

Corollary 2.8 Let L be a nilpotent Lie algebra and $V$ a finite dimensional indecomposable L-module. Then a basis can be chosen for $V$ with respect to which we obtain a matrix representation $\rho$ of $L$ of the form

$$
\rho(x)=\left(\begin{array}{ccccc}
\lambda(x) & & & & \\
& \cdot & & & * \\
\\
& & \cdot & & \\
& 0 & & & \\
& & & & \\
& & & \\
& & & & \\
& \text { for all } x \in L .
\end{array}\right) \quad \text {. }
$$

Proof. We can choose a basis as in Corollary 2.3 with respect to which each $\rho(x)$ is triangular. The generalised eigenspaces of $V$ with respect to $\rho(x)$ are all submodules of $V$ by Theorem 2.6 and $V$ is their direct sum. Since $V$ is indecomposable only one of the generalised eigenspaces is non-zero. Thus all the eigenvalues of $\rho(x)$ on $V$ are equal. Let this eigenvalue be $\lambda(x)$. Then the diagonal entries of the triangular matrix $\rho(x)$ are all equal to $\lambda(x)$.

We observe that the map $x \rightarrow \lambda(x)$ is a 1-dimensional representation of $L$, as it arises from a 1 -dimensional submodule of $V$.

We have seen from Proposition 2.5 and Theorem 2.6 how to obtain a direct decomposition of $V$ into submodules for any element $y \in L$. We may use this result to obtain a direct decomposition of $V$ into submodules which does not depend on the choice of any particular element of $L$.

Theorem 2.9 Let $L$ be a nilpotent Lie algebra and $V$ a finite dimensional $L$-module. For any 1-dimensional representation $\lambda$ of $L$ we define $V_{\lambda}=\{v \in V$; for each $x \in L$ there exists $N(x)$ such that $\left.(\rho(x)-\lambda(x) 1)^{N(x)} v=0\right\}$. Then

$$
V=\bigoplus_{\lambda} V_{\lambda}
$$

and each $V_{\lambda}$ is a submodule of $L$.

Proof. We first express $V$ as a direct sum of indecomposable $L$-modules. Each of these defines a 1-dimensional representation $\lambda$ of $L$ as in Corollary 2.8. Let $W_{\lambda}$ be the direct sum of all indecomposable components giving rise to $\lambda$. Then we have

$$
V=\bigoplus_{\lambda} W_{\lambda} .
$$

We shall show that $W_{\lambda}=V_{\lambda}$ and so that $W_{\lambda}$ is independent of the decomposition chosen into indecomposable components. It is clear that $W_{\lambda} \subset V_{\lambda}$ by Corollary 2.8. Suppose if possible that $W_{\lambda} \neq V_{\lambda}$. Then there exists $v \in V_{\lambda} \cap \bigoplus_{\mu \neq \lambda} W_{\mu}$ with $v \neq 0$. We write $v=\sum_{\mu \in S} w_{\mu}$ with $w_{\mu} \in W_{\mu}$, where the set $S$ is finite. Since $w_{\mu} \in W_{\mu}$ there exists $N_{\mu}$ such that $(\rho(x)-$ $\mu(x) 1)^{N_{\mu}} w_{\mu}=0$. Hence

$$
\prod_{\mu \in S}(\rho(x)-\mu(x) 1)^{N_{\mu}} v=0
$$

However, we also have $(\rho(x)-\lambda(x) 1)^{N_{\lambda}} v=0$.
We recall from Lemma 2.1 that the 1-dimensional representations of $L$ are in bijective correspondence with linear maps $L / L^{2} \rightarrow \mathbb{C}$. The vector space $L / L^{2}$ over $\mathbb{C}$ cannot be expressed as the union of finitely many proper subspaces. For each $\mu \in S$ the set of $x$ satisfying $\lambda(x)=\mu(x)$ is a proper subspace. Thus there exists $x \in L$ such that $\lambda(x) \neq \mu(x)$ for all $\mu \in S$. Thus the polynomials

$$
\prod_{\mu \in S}(t-\mu(x))^{N_{\mu}}, \quad(t-\lambda(x))^{N_{\lambda}}
$$

are coprime. Thus there exist polynomials $a(t), b(t) \in \mathbb{C}[t]$ such that

$$
a(t) \prod_{\mu \in S}(t-\mu(x))^{N_{\mu}}+b(t)(t-\lambda(x))^{N_{\lambda}}=1
$$

Hence

$$
a(\rho(x)) \prod_{\mu \in S}(\rho(x)-\mu(x) 1)^{N_{\mu}} v+b(\rho(x))(\rho(x)-\lambda(x) 1)^{N_{\lambda}} v=v
$$

The left-hand side of this expression is zero, as we have seen above. Thus $v=0$, a contradiction. Hence $V_{\lambda}=W_{\lambda}, V=\bigoplus_{\lambda} V_{\lambda}$ and each $V_{\lambda}$ is a submodule of $V$.

A 1-dimensional representation $\lambda$ of $L$ is called a weight of $V$ if $V_{\lambda} \neq 0$, and $V_{\lambda}$ is called the weight space of $\lambda$. The decomposition $V=\bigoplus_{\lambda} V_{\lambda}$ is called the weight space decomposition of $V$. It follows from Corollary 2.8 that a
basis can be chosen for $V_{\lambda}$ with respect to which the matrix representation of $L$ on $V_{\lambda}$ has form

$$
\rho(x)=\left(\begin{array}{ccccc}
\lambda(x) & & & & \\
& \cdot & & * & \\
& & \cdot & & \\
& & & \cdot & \\
& 0 & & & \\
& & & & \lambda(x)
\end{array}\right) \quad \text { for each } x \in L
$$

We shall make frequent use of the weight space decomposition in subsequent chapters.

We next prove a theorem of Engel which gives a useful characterisation of nilpotent Lie algebras in terms of the adjoint representation.

Theorem 2.10 (Engel's theorem). A Lie algebra L is nilpotent if and only if ad $x: L \rightarrow L$ is nilpotent for each $x \in L$.

Proof. Suppose $L$ is nilpotent. Then $L^{n}=O$ for some $n$. Let $y \in L$. Then we have

$$
\operatorname{ad} x \cdot y \in L^{2}, \quad(\operatorname{ad} x)^{2} \cdot y \in L^{3}, \quad \ldots
$$

and so $(\operatorname{ad} x)^{n-1} y=0$ for each $y \in \mathrm{~L}$. Thus $(\operatorname{ad} x)^{n-1}=0$ and so ad $x$ is a nilpotent linear map.

Now suppose conversely that ad $x$ is a nilpotent linear map for each $x \in L$. We wish to show $L$ is nilpotent. We suppose if possible that this is false and let $H$ be a maximal nilpotent subalgebra of $L$. Thus $H$ is nilpotent but any subalgebra properly containing $H$ is not nilpotent. We may regard $L$ as an $H$-module. Then $H$ is an $H$-submodule of $L$ and we can find an $H$-submodule $M$ of $L$ containing $H$ such that $M / H$ is an irreducible $H$-module. We have

$$
\operatorname{dim}(M / H)=1 \quad \text { by Theorem 2.2 }
$$

Moreover the 1-dimensional representation of $H$ afforded by $M / H$ is the zero representation, as otherwise ad $x$ would fail to be nilpotent for some $x \in H$. Hence we have $[H M] \subset H$. Now there exists $x \in M$ such that

$$
M=H \oplus \mathbb{C} x
$$

We have

$$
[M M] \subset[H H]+[H x] \subset H
$$

Thus $M$ is a subalgebra of $L$ and $H$ is an ideal of $M$.

We shall show that for each positive integer $i$ there exists a positive integer $e(i)$ such that

$$
M^{e(i)} \subset H^{i}
$$

This is true for $i=1$ since $M^{2} \subset H$. We prove it by induction on $i$. Assume that $M^{e(r)} \subset H^{r}$. Then

$$
M^{e(r)+1}=\left[M^{e(r)}, H+\mathbb{C} x\right] \subset H^{r+1}+\left[M^{e(r)}, x\right]
$$

Hence $M^{e(r)+1} \subset H^{r+1}+\operatorname{ad} x \cdot M^{e(r)}$.
We shall show that

$$
M^{e(r)+j} \subset H^{r+1}+(\operatorname{ad} x)^{j} \cdot M^{e(r)}
$$

for each positive integer $j$. This is true for $j=1$. Assuming it inductively for $j$ we have

$$
\begin{aligned}
M^{e(r)+j+1} & \subset\left[H^{r+1}+(\operatorname{ad} x)^{j} \cdot M^{e(r)}, M\right] \\
& \subset H^{r+1}+\left[(\operatorname{ad} x)^{j} M^{e(r)}, H+\mathbb{C} x\right] \\
& \subset H^{r+1}+(\operatorname{ad} x)^{j+1} M^{e(r)}
\end{aligned}
$$

since $H^{r+1}$ is an ideal of $M$ and $(\operatorname{ad} x)^{j} M^{e(r)} \subset H^{r}$. Thus we have shown

$$
M^{e(r)+j} \subset H^{r+1}+(\operatorname{ad} x)^{j} M^{e(r)} \quad \text { for all } j
$$

Now we know that $(\operatorname{ad} x)^{j}=0$ when $j$ is sufficiently large. For such $j$ we have

$$
M^{e(r)+j} \subset H^{r+1}
$$

Thus we define $e(r+1)=e(r)+j$ and then $M^{e(r+1)} \subset H^{r+1}$ as required.
Now $H$ is nilpotent so $H^{i}=O$ for $i$ sufficiently large. For such $i$ we have $M^{e(i)}=O$. Thus $M$ is nilpotent. But this contradicts the maximality of $H$. Thus our initial assumption was incorrect and so $L$ must be nilpotent.

Corollary 2.11 A Lie algebra L is nilpotent if and only if $L$ has a basis with respect to which the adjoint representation of $L$ has form

$$
\rho(x)=\left(\begin{array}{cccccc}
0 & & & & & \\
& 0 & & & & \\
& & \cdot & & * & \\
& & & \cdot & & \\
& 0 & & & & \\
& & & & 0 & \\
& & & & & 0
\end{array}\right) \quad \text { for all } x \in L
$$

Proof. Suppose $L$ is nilpotent. Then $L$ has a series of ideals

$$
L \supset L^{2} \supset L^{3} \supset \cdots \supset L^{r}=O \quad \text { for some } r
$$

We refine this series by choosing a sequence of subspaces between consecutive terms, each of codimension 1 in its predecessor. Such subspaces are automatically ideals of $L$ since if $L^{i} \supset I \supset L^{i+1}$ we have

$$
[I L] \subset\left[L^{i} L\right]=L^{i+1} \subset I
$$

Thus we have a chain of ideals

$$
L=I_{n} \supset I_{n-1} \supset \cdots \supset I_{1} \supset I_{0}=O
$$

with $\operatorname{dim} I_{k}=k$ and $\left[L I_{k}\right] \subset I_{k-1}$. By choosing a basis of $L$ adapted to this chain of ideals the map ad $x: L \rightarrow L$ is represented by a matrix $\rho(x)$ of zerotriangular form (i.e. triangular with zeros on the diagonal).

Conversely if $L$ has a basis with respect to which ad $x$ is represented by a zero-triangular matrix $\rho(x)$ for all $x \in L$, we have $\rho(x)$ nilpotent and so $\operatorname{ad} x$ is nilpotent. Thus $L$ must be a nilpotent Lie algebra by Engel's theorem (Theorem 2.10).

## 3

## Cartan subalgebras

### 3.1 Existence of Cartan subalgebras

Let $H$ be a subalgebra of a Lie algebra $L$. Let

$$
N(H)=\{x \in L ;[h x] \in H \quad \text { for all } h \in H\}
$$

$N(H)$ is called the normaliser of $H$.

Lemma 3.1 (i) $N(H)$ is a subalgebra of $L$.
(ii) $H$ is an ideal of $N(H)$.
(iii) $N(H)$ is the largest subalgebra of $L$ containing $H$ as an ideal.

Proof. (i) Let $x, y \in N(H)$. Then

$$
[h[x y]]=[[y h] x]+[[h x] y] \in H .
$$

Hence $[x y] \in N(H)$ and $N(H)$ is a subalgebra.
(ii) This is clear from the definition of $N(H)$.
(iii) If $H$ is an ideal of $M$ then $[H M] \subset H$ so $M \subset N(H)$.

Definition $A$ subalgebra $H$ of $L$ is called a Cartan subalgebra if $H$ is nilpotent and $H=N(H)$. Cartan subalgebras play a very important role in the theory of semisimple Lie algebras. Our aim in this section is to show that $L$ contains a Cartan subalgebra.

Let us take an element $x \in L$ and consider the linear map ad $x: L \rightarrow L$. Let $L_{0, x}$ be the generalised eigenspace of ad $x$ with eigenvalue 0 . Thus $L_{0, x}=$ $\left\{y \in L\right.$; there exists $n$ such that $\left.(\operatorname{ad} x)^{n} y=0\right\}$, and $L_{0, x}$ will be called the null component of $L$ with respect to $x$.

An element $x \in L$ is called regular if $\operatorname{dim} L_{0, x}$ is as small as possible. The Lie algebra $L$ certainly contains regular elements.

Theorem 3.2 Let $x$ be a regular element of $L$. Then the null component $L_{0, x}$ is a Cartan subalgebra of $L$.

Proof. Let $H=L_{0, x}$. We must show that $H$ is a subalgebra of $L$, that $H$ is nilpotent, and that $H=N(H)$.

We first show that $H$ is a subalgebra. Let $y, z \in H$. We must show that $[y z] \in H$. By Proposition 2.7 we have

$$
(\operatorname{ad} x)^{n}[y z]=\sum_{i=0}^{n}\binom{n}{i}\left[(\operatorname{ad} x)^{i} y,(\operatorname{ad} x)^{n-i} z\right]
$$

(We take $V=L, \alpha=\beta=0$ in Proposition 2.7 to obtain this.) Since $y \in H$ we have

$$
(\operatorname{ad} x)^{i} y=0 \quad \text { if } i \text { is sufficiently large. }
$$

Since $z \in H$

$$
(\operatorname{ad} x)^{n-i} z=0 \quad \text { if } n-i \text { is sufficiently large. }
$$

Hence $(\operatorname{ad} x)^{n}[y z]=0$ if $n$ is sufficiently large. Thus $[y z] \in H$ and $H$ is a subalgebra of $L$.

We next show that $H$ is nilpotent. To do this we shall prove that all the matrices in the adjoint representation of $H$ are nilpotent and use Engel's theorem (Theorem 2.10). Let $\operatorname{dim} H=l$ and $b_{1}, \ldots, b_{l}$ be a basis for $H$. Let

$$
y=\lambda_{1} b_{1}+\cdots+\lambda_{l} b_{l} \in H \quad \lambda_{1}, \ldots, \lambda_{l} \in \mathbb{C} .
$$

Consider the linear map ad $y: L \rightarrow L$. We have ad $y: H \rightarrow H$ since $H$ is a subalgebra and we obtain an induced map ad $y: L / H \rightarrow L / H$.

Let $\chi(t)$ be the characteristic polynomial of ad $y$ on $L, \chi_{1}(t)$ be its characteristic polynomial on $H$ and $\chi_{2}(t)$ be its characteristic polynomial on $L / H$. Then we have

$$
\chi(t)=\chi_{1}(t) \chi_{2}(t)
$$

Since $\chi(t)=\operatorname{det}(t 1-\operatorname{ad} y)$ and $y$ depends linearly on $\lambda_{1}, \ldots, \lambda_{l}$ we see that the coefficients of $\chi(t)$ are polynomial functions of $\lambda_{1}, \ldots, \lambda_{l}$. The same applies to $\chi_{1}(t)$ and $\chi_{2}(t)$. Let

$$
\chi_{2}(t)=d_{0}+d_{1} t+d_{2} t^{2}+\cdots
$$

where $d_{0}, d_{1}, d_{2}, \ldots$ are polynomial functions of $\lambda_{1}, \ldots, \lambda_{l}$. We claim that $d_{0}$ is not the zero polynomial. For in the special case when $y=x$ we know that
all eigenvalues of ad $y$ on $L / H$ are non-zero, so $\chi_{2}(t)$ has non-zero constant term. Let

$$
\chi_{1}(t)=t^{m}\left(c_{0}+c_{1} t+c_{2} t^{2}+\cdots\right)
$$

where $c_{0}, c_{1}, c_{2}, \ldots$ are polynomial functions of $\lambda_{1}, \ldots, \lambda_{l}$ and $c_{0}$ is not the zero polynomial. We have

$$
m \leq l=\operatorname{deg} \chi_{1}(t)
$$

We then have

$$
\chi(t)=t^{m}\left(c_{0} d_{0}+\text { terms involving positive powers of } t\right)
$$

Now $c_{0} d_{0}$ is not the zero polynomial so we can choose $\lambda_{1}, \ldots, \lambda_{l} \in \mathbb{C}$ to make $c_{0} d_{0}$ non-zero. For such an element $y \in H$ we have

$$
\operatorname{dim} L_{0, y}=m
$$

Since $x$ is regular and $\operatorname{dim} L_{0, x}=l$ we have $m \geq l$. Since we also know $m \leq l$ we have $m=l$. Now $\chi_{1}(t)$ has degree $l$ and is divisible by $t^{l}$, hence

$$
\chi_{1}(t)=t^{l} .
$$

It follows by the Cayley-Hamilton theorem that $(\operatorname{ad} y)^{l}: H \rightarrow H$ is zero. Hence by Engel's theorem we deduce that $H$ is nilpotent.

Finally we show that $H=N(H)$. It is certainly true that $H \subset N(H)$. So let $z \in N(H)$. Then $[x z] \in H$. Thus

$$
(\operatorname{ad} x)^{n}[x z]=0 \quad \text { for some } n
$$

But then $(\operatorname{ad} x)^{n+1} z=0$ and so $z \in H$. Thus $H=N(H)$ and we have shown that $H$ is a Cartan subalgebra of $L$.

### 3.2 Derivations and automorphisms

A derivation of a Lie algebra $L$ is a linear map $\delta: L \rightarrow L$ such that

$$
\delta[x y]=[\delta x, y]+[x, \delta y] \quad \text { for all } x, y \in L
$$

Lemma 3.3 Let $x \in L$. Then $\operatorname{ad} x$ is a derivation of $L$.

Proof. ad $x[y z]=[\operatorname{ad} x \cdot y, z]+[y, \operatorname{ad} x \cdot z]$ by the Jacobi identity.
An automorphism of $L$ is an isomorphism $\theta: L \rightarrow L$. The automorphisms of $L$ form a group Aut $L$ under composition.

Proposition 3.4 Let $\delta$ be a nilpotent derivation of L. Then $\exp \delta$ is an automorphism of $L$.

Proof. Since $\delta$ is nilpotent we have $\delta^{n}=0$ for some $n$. Then we have

$$
\exp \delta=\sum_{r=0}^{n-1} \frac{\delta^{r}}{r!}
$$

The map $\exp \delta: L \rightarrow L$ is clearly linear. Let $x, y \in L$. Then

$$
\begin{aligned}
& \delta[x y]=[\delta x, y]+[x, \delta y] \\
& \delta^{r}[x y]=\sum_{i=0}^{r}\binom{r}{i}\left[\delta^{i} x, \delta^{r-i} y\right]
\end{aligned}
$$

as is easily seen by induction on $r$. Hence

$$
\begin{aligned}
\exp \delta \cdot[x y] & =\sum_{r \geq 0} \sum_{i=0}^{r} \frac{1}{r!}\binom{r}{i}\left[\delta^{i} x, \delta^{r-i} y\right]=\sum_{i \geq 0} \sum_{j \geq 0} \frac{1}{i!j!}\left[\delta^{i} x, \delta^{j} y\right] \\
& =\left[\sum_{i \geq 0} \frac{1}{i!} \delta^{i} x, \sum_{j \geq 0} \frac{1}{j!} \delta^{j} y\right]=[\exp \delta \cdot x, \exp \delta \cdot y]
\end{aligned}
$$

Thus $\exp \delta: L \rightarrow L$ is a homomorphism. Similarly $\exp (-\delta)$ is a homomorphism and we have $\exp \delta \exp (-\delta)=1$. Thus $\exp \delta: L \rightarrow L$ is an automorphism.

The subgroup of Aut $L$ generated by all automorphisms exp ad $x$ for all $x \in L$ with ad $x$ nilpotent is called the group of inner automorphisms Inn $L$. Every element of Inn $L$ has form

$$
\exp \text { ad } x_{1} \cdot \exp \text { ad } x_{2} \cdots \cdot \exp \text { ad } x_{r}
$$

where $x_{1}, \ldots, x_{r} \in L$ and ad $x_{1}, \ldots$, ad $x_{r}$ are all nilpotent.

Lemma 3.5 Inn $L$ is a normal subgroup of Aut $L$.
Proof. Let $\theta \in$ Aut $L$. It is sufficient to show that $\theta(\exp$ ad $x) \theta^{-1} \in \operatorname{Inn} L$ for all $x \in L$ with ad $x$ nilpotent. Now we have

$$
\theta(\operatorname{ad} x) \theta^{-1} y=\theta\left[x, \theta^{-1} y\right]=[\theta x, y]=(\operatorname{ad} \theta x) \cdot y
$$

for all $y \in L$. Hence

$$
\theta(\operatorname{ad} x) \theta^{-1}=\operatorname{ad} \theta x
$$

It follows that

$$
\theta(\exp \operatorname{ad} x) \theta^{-1}=\exp \operatorname{ad}(\theta x) \in \operatorname{Inn} L
$$

Thus $\operatorname{Inn} L$ is normal in Aut $L$.
Two subalgebras $M_{1}, M_{2}$ of $L$ are called conjugate in $L$ if there exists $\theta \in$ Inn $L$ such that $\theta\left(M_{1}\right)=M_{2}$.

We wish to show that any two Cartan subalgebras of $L$ are conjugate in $L$. However, we first need some concepts from algebraic geometry.

### 3.3 Ideas from algebraic geometry

Let $H$ be a nilpotent subalgebra of a Lie algebra $L$ and regard $L$ as an $H$-module. Then we obtain a decomposition

$$
L=\bigoplus L_{\lambda}
$$

as in Theorem 2.9, where
$L_{\lambda}=\left\{x \in L ;\right.$ for each $h \in H$ there exists $n$ such that $\left.(\operatorname{ad} h-\lambda(h) 1)^{n} x=0\right\}$.

Now $H$ lies in $L_{0}$ by Corollary 2.11. We shall suppose that the nilpotent subalgebra $H$ satisfies the condition $H=L_{0}$. Then there exist 1-dimensional representations $\lambda_{1}, \ldots, \lambda_{r}$ of $H$ with $\lambda_{1} \neq 0, \ldots, \lambda_{r} \neq 0$ and

$$
L=H \oplus L_{\lambda_{1}} \oplus \cdots \oplus L_{\lambda_{r}} .
$$

Given $x \in L$ we then have

$$
x=x_{0}+x_{1}+\cdots+x_{r}
$$

with $x_{0} \in H$ and $x_{i} \in L_{\lambda_{i}}$ for $i=1, \ldots, r$. We claim that ad $x_{i}: L \rightarrow L$ is nilpotent when $i \neq 0$.

To see this let $\mu: H \rightarrow \mathbb{C}$ be a weight of the $H$-module $L$ and let $y \in L_{\mu}$. Then by Proposition 2.7 we have

$$
\begin{aligned}
\left(\operatorname{ad} h-\mu(h) 1-\lambda_{i}(h) 1\right)^{n}\left[x_{i} y\right]=\sum_{j=0}^{n}\binom{n}{j} & {\left[\left(\operatorname{ad} h-\lambda_{i}(h) 1\right)^{j} x_{i}\right.} \\
& \left.(\operatorname{ad} h-\mu(h) 1)^{n-j} y\right] .
\end{aligned}
$$

Because $x_{i} \in L_{\lambda_{i}}$ then $\left(\operatorname{ad} h-\lambda_{i}(h) 1\right)^{j} x_{i}=0$ if $j$ is sufficiently large. Since $y \in L_{\mu}$ then $(\operatorname{ad} h-\mu(h) 1)^{n-j} y=0$ if $n-j$ is sufficiently large. Thus

$$
\left(\operatorname{ad} h-\mu(h) 1-\lambda_{i}(h) 1\right)^{n}\left[x_{i} y\right]=0
$$

if $n$ is sufficiently large, and so $\left[x_{i} y\right] \in L_{\lambda_{i}+\mu}$. Thus we have

$$
\operatorname{ad} x_{i} \cdot L_{\mu} \subset L_{\lambda_{i}+\mu}
$$

Since $\lambda_{i} \neq 0$ and there are only finitely many $\mu: H \rightarrow \mathbb{C}$ for which $L_{\mu} \neq 0$ we see that $\left(\operatorname{ad} x_{i}\right)^{N}=0$ if $N$ is sufficiently large. Thus ad $x_{i}$ is nilpotent.

We deduce that $\exp$ ad $x_{i} \in$ Aut $L$ for $i \neq 0$. We now define a map $f: L \rightarrow L$ by

$$
f(x)=\exp \text { ad } x_{1} \cdot \exp \text { ad } x_{2} \cdots \cdots \exp \text { ad } x_{r} \cdot x_{0}
$$

We shall discuss some properties of this function $f$. We choose a basis $\left\{b_{i j}\right\}$ of $L$ for $0 \leq i \leq r$ where for fixed $i$ the elements $b_{i j}$ form a basis of $L_{\lambda_{i}}$ with respect to which the elements of $H$ are represented by triangular matrices, as in Corollary 2.3. Here $\lambda_{0}=0$.

Lemma $3.6 f: L \rightarrow L$ is a polynomial function. Thus

$$
f\left(\sum \lambda_{i j} b_{i j}\right)=\sum \mu_{i j} b_{i j}
$$

where each $\mu_{i j}$ is a polynomial in the $\lambda_{k l}$.
Proof. Each map ad $x_{i}: L \rightarrow L$ is linear. Also we have

$$
\exp \operatorname{ad} x_{i}=\sum_{k=0}^{N} \frac{\left(\operatorname{ad} x_{i}\right)^{k}}{k!} \quad \text { for some } N
$$

since ad $x_{i}$ is nilpotent. Thus $\exp$ ad $x_{i}: L \rightarrow L$ is a polynomial function.
The given map $f$ is a composition of the linear map $x \rightarrow x_{0}$ with polynomial functions exp ad $x_{i}$ for $i>0$, so is a polynomial function.

We write $\mu_{i j}=f_{i j}\left(\lambda_{k l}\right)$ where $f_{i j}$ is a polynomial. We define the Jacobian matrix

$$
J(f)=\left(\partial f_{i j} / \partial \lambda_{k l}\right)
$$

and the Jacobian determinant $\operatorname{det} J(f)$ of $f . \operatorname{det} J(f)$ is an element of the polynomial ring $\mathbb{C}\left[\lambda_{k l}\right]$.

Proposition $3.7 \operatorname{det} J(f)$ is not the zero polynomial.

Proof. We shall show det $J(f)$ is not the zero polynomial by showing that it is non-zero when evaluated at a carefully chosen element of $H$. So let $h \in H$ and consider $\left(\partial f_{i j} / \partial \lambda_{k l}\right)_{h}$.

First suppose $k \neq 0$. Then

$$
\begin{aligned}
\left(\partial f / \partial \lambda_{k l}\right)_{h} & =\lim _{t \rightarrow 0} \frac{f\left(h+t b_{k l}\right)-f(h)}{t} \\
& =\lim _{t \rightarrow 0} \frac{\left(\exp \operatorname{ad} t b_{k l}\right) h-h}{t} \\
& =\lim _{t \rightarrow 0} \frac{h+t\left[b_{k l}, h\right]+\cdots-h}{t} \\
& =\left[b_{k l}, h\right]=-\left[h b_{k l}\right] \\
& =-\lambda_{k}(h) b_{k l}+\text { a linear combination of } b_{k 1}, \ldots, b_{k l-1}
\end{aligned}
$$

Next suppose $k=0$. Then

$$
\begin{aligned}
\left(\partial f / \partial \lambda_{0 l}\right)_{h} & =\lim _{t \rightarrow 0} \frac{f\left(h+t b_{0 l}\right)-f(h)}{t} \\
& =\lim _{t \rightarrow 0} \frac{h+t b_{0 l}-h}{t}=b_{0 l} .
\end{aligned}
$$

Thus $J(f)_{h}$ is a block matrix of form
and so $(\operatorname{det} J(f))_{h}= \pm \prod_{i=1}^{r} \lambda_{i}(h)^{d_{i}}$ where $d_{i}=\operatorname{dim} L_{\lambda_{i}}$.
Now the linear maps $\lambda_{i}: H \rightarrow \mathbb{C}$ for $i=1, \ldots, v$ are all non-zero. Thus we can find an element $h \in H$ with $\lambda_{i}(h) \neq 0$ for $i=1, \ldots, v$. For such an element $h$ we have $(\operatorname{det} J(f))_{h} \neq 0$. Hence $\operatorname{det} J(f)$ is not the zero polynomial.

Proposition 3.8 The polynomial functions $f_{i j}$ are algebraically independent.

Proof. Suppose if possible that there is a non-zero polynomial $F\left(x_{i j}\right) \in \mathbb{C}\left[x_{i j}\right]$ such that $F\left(f_{i j}\right)=0$. We choose such a polynomial $F$ whose total degree in the variables $x_{i j}$ is as small as possible. Then

$$
\frac{\partial}{\partial \lambda_{k l}} F\left(f_{i j}\right)=0
$$

and so

$$
\sum_{i, j} \frac{\partial F}{\partial f_{i j}} \frac{\partial f_{i j}}{\partial \lambda_{k l}}=0
$$

Let $v$ be the vector $\left(\partial F / \partial f_{i j}\right)$. Then

$$
v J(f)=(0, \ldots, 0)
$$

Since $\operatorname{det} J(f)$ is non-zero this implies that $v=(0, \ldots, 0)$, that is

$$
\partial F / \partial f_{i j}=0 \quad \text { for each } f_{i j}
$$

Now $\partial F / \partial x_{i j}$ is a polynomial in $\mathbb{C}\left[x_{i j}\right]$ of smaller total degree than $F$. By the choice of $F \partial F / \partial x_{i j}$ must be the zero polynomial. Hence $F$ does not involve the variable $x_{i j}$. Since this is true for all $x_{i j} F$ must be a constant. Since $F\left(f_{i j}\right)=0$ this constant must be zero. Thus $F$ is the zero polynomial and we have a contradiction.

Let $B=\mathbb{C}\left[f_{i j}\right]$ be the polynomial ring in the $f_{i j}$ and $A=\mathbb{C}\left[\lambda_{i j}\right]$ the polynomial ring in the $\lambda_{i j}$. We have a homomorphism $\theta: B \rightarrow A$ uniquely determined by

$$
\theta\left(f_{i j}\right)=f_{i j}\left(\lambda_{k l}\right) \in A
$$

Proposition 3.9 The homomorphism $\theta: B \rightarrow A$ is injective.

Proof. Suppose $F \in B$ satisfies $\theta(F)=0$. Then $F\left(f_{i j}\right)=0$, regarded as a function of the $\lambda_{k l}$. Since the $f_{i j}$ are algebraically independent this implies that $F=0$. Thus $\theta$ is injective.

Thus we may regard $B$ as a subring of $A . A$ and $B$ are integral domains with a common identity element and $A$ is finitely generated over $B$. We next prove a general result which applies to this situation.

Proposition 3.10 Let $A$ and $B$ be integral domains such that $B \subset A, A, B$ have a common identity element 1 , and $A$ is finitely generated over $B$. Let $p$ be a non-zero element of $A$. Then there exists a non-zero element $q$ of $B$ such that any homomorphism $\phi: B \rightarrow \mathbb{C}$ with $\phi(q) \neq 0$ can be extended to a homomorphism $\psi: A \rightarrow \mathbb{C}$ with $\psi(p) \neq 0$.

Proof. We may assume that $A$ is generated over $B$ by a single element $\zeta$. For then by iterating the process we can prove the result when $A$ is finitely generated over $B$. Thus we assume that $A=B[\zeta]$ for some $\zeta \in A$.

Suppose first that $\zeta$ is transcendental over $B$. Given a non-zero element $p=p(\zeta) \in A$ we choose $q \in B$ to be one of the non-zero coefficients of $p(\zeta)$. Suppose we are given a homomorphism $\phi: B \rightarrow \mathbb{C}$ with $\phi(q) \neq 0$. We write $\phi(b)=\bar{b} \in \mathbb{C}$. By applying $\phi$ to the coefficients of $p(\zeta)$ we obtain $\bar{p}(\zeta) \in \mathbb{C}[\zeta]$. The element $\bar{p}(\zeta)$ is not the zero polynomial since $\phi(q) \neq 0$. We can find an element $\beta \in \mathbb{C}$ with $\bar{p}(\beta) \neq 0$. We now define a homomorphism $\psi: A \rightarrow \mathbb{C}$ by

$$
\psi(g(\zeta))=\bar{g}(\beta)
$$

$\psi$ is well defined since $\zeta$ is transcendental over $B$, and $\psi$ is a homomorphism, being a composite of the homomorphisms

$$
\begin{align*}
& A=B[\zeta] \rightarrow \mathbb{C}[\zeta] \rightarrow \mathbb{C}  \tag{3.1}\\
& g(\zeta) \quad \rightarrow \bar{g}(\zeta) \rightarrow \bar{g}(\beta) \tag{3.2}
\end{align*}
$$

$\psi$ clearly extends $\phi$. Finally we have $\psi(p)=\bar{p}(\beta) \neq 0$.
Next suppose that $\zeta$ is algebraic over $B$. Then we can find $f(t) \in B[t]$ of minimal degree such that $f(\zeta)=0$. We write

$$
f(t)=b_{0} t^{n}+b_{1} t^{n-1}+\cdots+b_{n} \quad b_{0} \neq 0
$$

Now let $g(t)$ be any polynomial in $B[t]$ satisfying $g(\zeta)=0$. We divide $g(t)$ by $f(t)$ using the Euclidean algorithm. We are working over an integral domain $B$ rather than over a field. However, provided we multiply $g(t)$ by a sufficiently high power of the leading coefficient $b_{0}$ of $f(t)$ we can carry out the Euclidean process over $B$. We thus obtain

$$
b_{0}^{k} g(t)=u(t) f(t)+v(t)
$$

where $u(t), v(t) \in B[t]$ and $\operatorname{deg} v(t)<\operatorname{deg} f(t)$. Thus

$$
v(\zeta)=b_{0}^{k} g(\zeta)-u(\zeta) f(\zeta)=0
$$

Since $\operatorname{deg} v(t)<\operatorname{deg} f(t)$ this implies that $v(t)=0$. Hence

$$
b_{0}^{k} g(t)=u(t) f(t)
$$

Let $p$ be the given non-zero element of $A$. The element $p$ is algebraic over $B$ since $A$ is generated over $B$ by the single algebraic element $\zeta$. Thus there exists a polynomial $h(t) \in B[t]$ with non-zero constant term $h_{m}$ such that $h(p)=0$. We define the element $q \in B$ by $q=b_{0} h_{m}$. Thus $q \neq 0$. We assume we are given a homomorphism $\phi: B \rightarrow \mathbb{C}$ with $\phi(q) \neq 0$. Then $\phi\left(b_{0}\right) \neq 0$ and
$\phi\left(h_{m}\right) \neq 0$. We write $\phi(b)=\bar{b} \in \mathbb{C}$. The polynomial $f(t) \in B[t]$ gives rise to a polynomial $\bar{f}(t) \in \mathbb{C}[t]$. We choose an element $\beta \in \mathbb{C}$ with $\bar{f}(\beta)=0$. We note that

$$
\bar{b}_{0}^{k} \bar{g}(\beta)=\bar{u}(\beta) \bar{f}(\beta)=0
$$

hence $\bar{g}(\beta)=0$ since $\bar{b}_{0}=\phi\left(b_{0}\right) \neq 0$.
We now define a homomorphism

$$
\psi: A \rightarrow \mathbb{C}
$$

by $\psi(g(\zeta))=\bar{g}(\beta)$. We note that the map $\psi$ is well defined, since we have shown that $g(\zeta)=0$ implies $\bar{g}(\beta)=0$. The map $\psi$ is a homomorphism since the maps

$$
\begin{aligned}
& B[t] \rightarrow \mathbb{C}[t] \rightarrow \mathbb{C} \\
& g(t) \rightarrow \bar{g}(t) \rightarrow \bar{g}(\beta)
\end{aligned}
$$

are homomorphisms. The definition of $\psi$ shows that $\psi$ extends $\phi$. Finally we show $\psi(p) \neq 0$. Since $h(p)=0$ we have $\bar{h}(\psi(p))=0$. However, the constant term of $h(t)$ is $\phi\left(h_{m}\right)$, which is non-zero. Since $\bar{h}(t)$ has non-zero constant term and $h(\psi(p))=0$ we must have $\psi(p) \neq 0$.

We now apply this result to our earlier situation. Let $d=\operatorname{dim} L$ and

$$
f: \mathbb{C}^{d} \rightarrow \mathbb{C}^{d}
$$

be the polynomial function

$$
\left(\lambda_{i j}\right) \rightarrow\left(f_{i j}\left(\lambda_{k l}\right)\right)
$$

We write $V=\mathbb{C}^{d}$ and for each polynomial $p \in \mathbb{C}\left[x_{i j}\right]$ we write

$$
V_{p}=\{v \in V ; p(v) \neq 0\} .
$$

Corollary 3.11 For each non-zero polynomial $p \in \mathbb{C}\left[x_{i j}\right]$ there exists a nonzero polynomial $q \in \mathbb{C}\left[x_{i j}\right]$ such that $f\left(V_{p}\right) \supset V_{q}$.

Proof. We apply Proposition 3.10 to the integral domains $B \subset A$ discussed earlier. Thus $A$ is the polynomial ring $\mathbb{C}\left[\lambda_{i j}\right]$ and $B$ is the polynomial ring $\mathbb{C}\left[f_{i j}\right]$. We choose a non-zero polynomial $p \in A$. Then there exists a non-zero polynomial $q \in B$ such that any homomorphism $\phi: B \rightarrow \mathbb{C}$ with $\phi(q) \neq 0$ can be extended to a homomorphism $\psi: A \rightarrow \mathbb{C}$ with $\psi(p) \neq 0$. This means that given any $v \in V_{q}$ we have $v=f(w)$ for some $w \in V_{p}$. Hence $V_{q} \subset f\left(V_{p}\right)$ as required.

### 3.4 Conjugacy of Cartan subalgebras

We showed in Theorem 3.2 that the null component $L_{0, x}$ of a regular element $x \in L$ is a Cartan subalgebra of $L$. We shall now show conversely that any Cartan subalgebra is the null component of some regular element. We shall then prove that, given two regular elements, their null components are conjugate in $L$.

Proposition 3.12 Let $H$ be a Cartan subalgebra of $L$. Then there exists a regular element $x \in L$ such that $H=L_{0, x}$.

Proof. Since $H$ is nilpotent we may regard $L$ as an $H$-module and decompose $L$ into weight spaces with respect to $H$ as in Theorem 2.9. $H$ lies in the zero weight space $L_{0}$ by Corollary 2.11. Since $H=N(H)$ we can show that $H=L_{0}$. For if $H \neq L_{0}$ the $H$-module $L_{0} / H$ will have a 1 -dimensional submodule $M / H$ on which $H$ acts with weight 0 . Hence $[H M] \subset H$ and so $M \subset N(H)$. This contradicts $H=N(H)$. Thus we have $H=L_{0}$. Let

$$
L=H \oplus L_{\lambda_{1}} \oplus \cdots \oplus L_{\lambda_{r}} \quad \lambda_{1}, \ldots, \lambda_{r} \neq 0
$$

be the weight space decomposition of $L$ with respect to $H$. Let $x \in L$ and

$$
x=x_{0}+x_{1}+\cdots+x_{r}
$$

with $x_{0} \in H$ and $x_{i} \in L_{\lambda_{i}}$ for $i \neq 0$. Then we can define a polynomial function $f: L \rightarrow L$ as in Section 3.3 with

$$
f(x)=\exp \text { ad } x_{1} \cdot \exp \text { ad } x_{2} \cdots \cdot \exp \text { ad } x_{r} \cdot x_{0}
$$

We define $p: L \rightarrow \mathbb{C}$ by

$$
p(x)=\lambda_{1}\left(x_{0}\right) \lambda_{2}\left(x_{0}\right) \cdots \lambda_{r}\left(x_{0}\right)
$$

Then $p$ is a polynomial function on $L . p$ is not the zero polynomial since we can find $x_{0} \in H$ for which each $\lambda_{i}\left(x_{0}\right) \neq 0$ for $i=1, \ldots, r$. Hence by Corollary 3.11 there exists a non-zero polynomial function $q: L \rightarrow \mathbb{C}$ such that $f\left(L_{p}\right) \supset L_{q}$.

We next consider the set $R$ of regular elements of $L$. Let $y \in L$ and

$$
\chi(y)=\operatorname{det}(t 1-\operatorname{ad} y)=t^{n}+\mu_{1}(y) t^{n-1}+\cdots+\mu_{n}(y)
$$

be the characteristic polynomial of ad $y$ on $L$. Then $\mu_{1}, \mu_{2}, \ldots, \mu_{n}$ are polynomial functions on $L$. There exists a unique integer $k$ such that $\mu_{n-k}$ is not the zero polynomial but $\mu_{n-k+1}, \ldots, \mu_{n}$ are identically zero. The generalised eigenspace of ad $y$ with eigenvalue 0 has dimension $k$ if $\mu_{n-k}(y) \neq 0$
and dimension greater than $k$ if $\mu_{n-k}(y)=0$. Thus $y$ is regular if and only if $\mu_{n-k}(y) \neq 0$.

Now there exists $y \in L$ such that $y \in L_{q} \cap R$. For we may choose $y$ with $\left(q \mu_{n-k}\right)(y) \neq 0$. Since $L_{q} \subset f\left(L_{p}\right)$ we can find $x \in L_{p}$ such that $f(x)=y$. Thus we have

$$
\exp \text { ad } x_{1} \cdot \exp \text { ad } x_{2} \cdots \exp \text { ad } x_{r} \cdot x_{0}=y
$$

Hence $x_{0}, y$ are conjugate elements of $L$. Since $y$ is regular, $x_{0}$ must also be regular. Since $x \in L_{p}$ we have

$$
\lambda_{1}\left(x_{0}\right) \lambda_{2}\left(x_{0}\right) \cdots \lambda_{r}\left(x_{0}\right) \neq 0
$$

Now $x_{0} \in H$ and $H$ is nilpotent, hence $L_{0, x_{0}} \supset H$ by Corollary 2.11. On the other hand $L_{0, x_{0}}$ cannot be larger than $H$ since

$$
\lambda_{1}\left(x_{0}\right) \neq 0, \ldots, \lambda_{r}\left(x_{0}\right) \neq 0
$$

Hence $H=L_{0, x_{0}}$ where $x_{0}$ is regular.
Theorem 3.13 Any two Cartan subalgebras of $L$ are conjugate.
Proof. Let $H, H^{\prime}$ be Cartan subalgebras of $L$. We regard $L$ as an $H$-module and decompose $L$ into weight spaces with respect to $H$. We have seen in the proof of Proposition 3.12 that $H=L_{0}$. Let the weight space decomposition be

$$
L=H \oplus L_{\lambda_{1}} \oplus \cdots \oplus L_{\lambda_{r}} \quad \lambda_{1}, \ldots, \lambda_{r} \neq 0
$$

For each $x \in L$ we have

$$
x=x_{0}+x_{1}+\cdots+x_{r}
$$

with $x_{0} \in H$ and $x_{i} \in L_{\lambda_{i}}$ for $i \neq 0$.
Now for each $x_{0} \in H$ we have $L_{0, x_{0}} \supset H$ and for some $x_{0} \in H$ we have $L_{0, x_{0}}=H$ since $H$ is a Cartan subalgebra. An element $x_{0} \in H$ is regular if and only if $L_{0, x_{0}}=H$. This is equivalent to the condition

$$
\lambda_{1}\left(x_{0}\right) \lambda_{2}\left(x_{0}\right) \cdots \lambda_{r}\left(x_{0}\right) \neq 0
$$

We now consider the polynomial function $f: L \rightarrow L$ defined by

$$
f(x)=\exp \text { ad } x_{1} \cdot \exp \text { ad } x_{2} \cdots \cdots \exp \text { ad } x_{r} \cdot x_{0}
$$

Let $p: L \rightarrow \mathbb{C}$ be the function given by

$$
p(x)=\lambda_{1}\left(x_{0}\right) \lambda_{2}\left(x_{0}\right) \cdots \lambda_{r}\left(x_{0}\right)
$$

where $p$ is a polynomial function on $L$ which is not identically zero, since $p(x)$ is non-zero when $x_{0}$ is a regular element of $H$. By Corollary 3.11 there exists a non-zero polynomial function $q: L \rightarrow \mathbb{C}$ such that $f\left(L_{p}\right) \supset L_{q}$.

We now start with the second Cartan subalgebra $H^{\prime}$. We can define a corresponding function $f^{\prime}: L \rightarrow L$ and a corresponding function $p^{\prime}: L \rightarrow \mathbb{C}$. There exists a non-zero polynomial function $q^{\prime}: L \rightarrow \mathbb{C}$ such that $f^{\prime}\left(L_{p^{\prime}}\right) \supset L_{q^{\prime}}$.

Now $L_{q} \cap L_{q^{\prime}}=\left\{x \in L ;\left(q q^{\prime}\right)(x) \neq 0\right\}$. Thus $L_{q} \cap L_{q^{\prime}}$ is non-empty. We choose $z \in L_{q} \cap L_{q^{\prime}}$. Thus $z \in f\left(L_{p}\right) \cap f^{\prime}\left(L_{p^{\prime}}\right)$. Thus there exists $x \in L$ with $z=f(x)$ and $p(x) \neq 0$. Similarly there exists $x^{\prime} \in L$ with $z=f^{\prime}\left(x^{\prime}\right)$ and $p^{\prime}(x) \neq 0$. Thus

$$
z=\exp \text { ad } x_{1} \cdot \exp \text { ad } x_{2} \cdots \cdots \exp \text { ad } x_{r} \cdot x_{0}
$$

and so $z$ is conjugate to $x_{0}$. Since $p(x) \neq 0 x_{0}$ is regular. Similarly $z$ is conjugate to $x_{0}^{\prime}$ and $x_{0}^{\prime}$ is regular. Thus we have found regular elements $x_{0} \in H$ and $x_{0}^{\prime} \in H^{\prime}$ such that $x_{0}, x_{0}^{\prime}$ are conjugate in $L$.

Now we have $H=L_{0, x_{0}}$ and $H^{\prime}=L_{0, x_{0}^{\prime}}$ since $x_{0}, x_{0}^{\prime}$ are regular. Thus an inner automorphism of $L$ which transforms $x_{0}$ to $x_{0}^{\prime}$ will transform $H$ to $H^{\prime}$. Hence $H, H^{\prime}$ are conjugate in $L$.

The dimension of the Cartan subalgebras of $L$ will be called the rank of $L$.

## 4

## The Cartan decomposition

### 4.1 Some properties of root spaces

Let $L$ be a Lie algebra and $H$ be a Cartan subalgebra of $L$. We regard $L$ as an $H$-module. Since $H$ is nilpotent we have a weight space decomposition

$$
L=\bigoplus_{\lambda} L_{\lambda}
$$

as in Theorem 2.9, where
$L_{\lambda}=\left\{x \in L\right.$; for each $h \in H$ there exists $n$ such that $\left.(\operatorname{ad} h-\lambda(h) 1)^{n} x=0\right\}$.

Proposition 4.1 $L_{0}=H$.
Proof. The algebra $H$ is contained in $L_{0}$ by Corollary 2.11. Suppose if possible that $H \neq L_{0}$. Then $L_{0} / H$ is an $H$-module, and this module contains a 1-dimensional submodule $M / H$ on which $H$ acts with weight 0 . Hence $[H M] \subset H$ and so $M \subset N(H)$. This implies $H \neq N(H)$, a contradiction.

The 1 -dimensional representations $\lambda$ of $H$ such that $\lambda \neq 0$ and $L_{\lambda} \neq O$ are called the roots of $L$ with respect to $H$. The set of roots of $L$ with respect to $H$ will be denoted by $\Phi$. Thus we have

$$
L=H \oplus\left(\bigoplus_{\alpha \in \Phi} L_{\alpha}\right)
$$

This decomposition is called the Cartan decomposition of $L$ with respect to $H$. $L_{\alpha}$ is called the root space of $\alpha$.

Proposition 4.2 Let $\lambda, \mu$ be 1-dimensional representations of $H$. Then

$$
\left[L_{\lambda}, L_{\mu}\right] \subset L_{\lambda+\mu} .
$$

Proof. Let $y \in L_{\lambda}, z \in L_{\mu}$. We show that $[y z] \in L_{\lambda+\mu}$. Let $x \in H$. Then by Proposition 2.7 we have
$(\operatorname{ad} x-\lambda(x) 1-\mu(x) 1)^{n}[y z]=\sum_{i=0}^{n}\binom{n}{i}\left[(\operatorname{ad} x-\lambda(x) 1)^{i} y,(\operatorname{ad} x-\mu(x) 1)^{n-i} z\right]$.
Since $y \in L_{\lambda} \quad(\operatorname{ad} x-\lambda(x) 1)^{i} y=0$ if $i$ is sufficiently large. Since $z \in L_{\mu}$ $(\operatorname{ad} x-\mu(x) 1)^{n-i} z=0$ if $n-i$ is sufficiently large. Hence

$$
(\operatorname{ad} x-\lambda(x) 1-\mu(x) 1)^{n}[y z]=0
$$

if $n$ is sufficiently large. This shows that $[y z] \in L_{\lambda+\mu}$.

Corollary 4.3 Let $\alpha, \beta \in \Phi$ be roots of $L$ with respect to $H$. Then

$$
\begin{array}{ll}
{\left[L_{\alpha}, L_{\beta}\right] \subset L_{\alpha+\beta}} & \text { if } \alpha+\beta \in \Phi \\
{\left[L_{\alpha}, L_{\beta}\right] \subset H} & \text { if } \beta=-\alpha \\
{\left[L_{\alpha}, L_{\beta}\right]=0} & \text { if } \alpha+\beta \neq 0 \text { and } \alpha+\beta \notin \Phi .
\end{array}
$$

Proof. This follows from Proposition 4.2 and the fact that $L_{0}=H$.

Proposition 4.4 Let $\alpha \in \Phi$ and consider the subspace $\left[L_{\alpha} L_{-\alpha}\right]$ of $H$. Given any $\beta \in \Phi$ there exists a number $r \in \mathbb{Q}$, depending on $\alpha$ and $\beta$, such that $\beta=r \alpha$ on $\left[L_{\alpha} L_{-\alpha}\right]$.

Proof. If $-\alpha$ is not a weight of $L$ with respect to $H$ then $L_{-\alpha}=O$ and there is nothing to prove. Thus we assume $-\alpha$ is a weight. Then $-\alpha \in \Phi$ since $\alpha \neq 0$.

We consider the functions $i \alpha+\beta: H \rightarrow \mathbb{C}$ where $i \in \mathbb{Z}$. Since $\Phi$ is finite there exist $p, q \in \mathbb{Z}$ with $p \geq 0, q \geq 0$ such that

$$
-p \alpha+\beta, \ldots, \beta, \ldots, q \alpha+\beta
$$

are all in $\Phi$ but $-(p+1) \alpha+\beta,(q+1) \alpha+\beta$ are not in $\Phi$. If either $-(p+1) \alpha+\beta=0$ or $(q+1) \alpha+\beta=0$ the result is obvious. Thus we assume $-(p+1) \alpha+\beta \neq 0,(q+1) \alpha+\beta \neq 0$. Thus $-(p+1) \alpha+\beta,(q+1) \alpha+\beta$ are not weights of $L$ with respect to $H$.

Let $M$ be the subspace of $L$ given by

$$
M=L_{-p \alpha+\beta} \oplus \cdots \oplus L_{q \alpha+\beta}
$$

Let $y \in L_{\alpha}, z \in L_{-\alpha}$. Let $x=[y z] \in\left[L_{\alpha} L_{-\alpha}\right]$. Then we have

$$
\begin{array}{ll}
\text { ad } y(M) \subset M & \text { by Proposition 4.2, since } L_{(q+1) \alpha+\beta}=O \\
\operatorname{ad} z(M) \subset M & \text { by Proposition 4.2, since } L_{-(p+1) \alpha+\beta}=O
\end{array}
$$

Thus

$$
\operatorname{ad} x(M)=(\operatorname{ad} y \operatorname{ad} z-\operatorname{ad} z \operatorname{ad} y) M \subset M
$$

We calculate the $\operatorname{trace} \operatorname{tr}_{M}$ ad $x$. Since $x \in H$ each weight space $L_{i \alpha+\beta}$ is invariant under ad $x$. Thus

$$
\operatorname{tr}_{M} \operatorname{ad} x=\sum_{i=-p}^{q} \operatorname{tr}_{L_{i \alpha+\beta}} \operatorname{ad} x
$$

Now ad $x$ acts on the weight space $L_{i \alpha+\beta}$ by means of a matrix of form

$$
\left(\begin{array}{cccc}
(i \alpha+\beta) x & & & * \\
& \cdot & & \\
& \cdot & \\
& & \cdot & \\
0 & & & (i \alpha+\beta) x
\end{array}\right)
$$

Thus $\operatorname{tr}_{L_{i \alpha+\beta}}$ ad $x=\operatorname{dim} L_{i \alpha+\beta} .(i \alpha+\beta)(x)$. Thus

$$
\begin{aligned}
\operatorname{tr}_{M} \text { ad } x & =\sum_{i=-p}^{q} \operatorname{dim} L_{i \alpha+\beta}(i \alpha(x)+\beta(x)) \\
& =\left(\sum_{i=-p}^{q} i \operatorname{dim} L_{i \alpha+\beta}\right) \alpha(x)+\left(\sum_{i=-p}^{q} \operatorname{dim} L_{i \alpha+\beta}\right) \beta(x) .
\end{aligned}
$$

On the other hand we have

$$
\begin{aligned}
\operatorname{tr}_{M} \text { ad } x & =\operatorname{tr}_{M}(\operatorname{ad} y \operatorname{ad} z-\operatorname{ad} z \operatorname{ad} y) \\
& =\operatorname{tr}_{M}(\operatorname{ad} y \operatorname{ad} z)-\operatorname{tr}_{M}(\operatorname{ad} z \operatorname{ad} y)=0 .
\end{aligned}
$$

Hence

$$
\left(\sum_{i=-p}^{q} i \operatorname{dim} L_{i \alpha+\beta}\right) \alpha(x)+\left(\sum_{i=-p}^{q} \operatorname{dim} L_{i \alpha+\beta}\right) \beta(x)=0 .
$$

Moreover $\operatorname{dim} L_{i \alpha+\beta}>0$ for $-p \leq i \leq q$. Hence for $x \in\left[L_{\alpha} L_{-\alpha}\right]$ we have

$$
\beta(x)=-\frac{\left(\sum_{i=-p}^{q} i \operatorname{dim} L_{i \alpha+\beta}\right)}{\left(\sum_{i=-p}^{q} \operatorname{dim} L_{i \alpha+\beta}\right)} \alpha(x) .
$$

Thus $\beta(x)=r \alpha(x)$ for some $r \in \mathbb{Q}$ independent of $x$. Hence $\beta=r \alpha$ on $\left[L_{\alpha} L_{-\alpha}\right]$.

### 4.2 The Killing form

In order to make further progress in understanding the Cartan decomposition of $L$ we introduce a bilinear form on $L$ called the Killing form. We define a map

$$
\begin{aligned}
& L \times L \rightarrow \mathbb{C} \\
& x, y \rightarrow\langle x, y\rangle
\end{aligned}
$$

given by

$$
\langle x, y\rangle=\operatorname{tr}(\operatorname{ad} x \text { ad } y)
$$

We have ad $x: L \rightarrow L$, ad $y: L \rightarrow L$ and ad $x$ ad $y: L \rightarrow L$, so $\operatorname{tr}(\operatorname{ad} x \operatorname{ad} y) \in \mathbb{C}$.

Proposition 4.5 (i) $\langle x, y\rangle$ is bilinear, i.e. linear in $x$ and $y$.
(ii) $\langle x, y\rangle$ is symmetric, i.e. $\langle y, x\rangle=\langle x, y\rangle$.
(iii) $\langle x, y\rangle$ is invariant, i.e.

$$
\langle[x y], z\rangle=\langle x,[y z]\rangle \quad \text { for all } x, y, z \in L
$$

Proof. (i) is clear from the definition.
(ii) follows from the fact that $\operatorname{tr} A B=\operatorname{tr} B A$.
(iii) $\langle[x y], z\rangle=\operatorname{tr}(\operatorname{ad}[x y] \operatorname{ad} z)=\operatorname{tr}((\operatorname{ad} x$ ad $y-\operatorname{ad} y \operatorname{ad} x) \operatorname{ad} z)$

$$
\begin{aligned}
& =\operatorname{tr}(\operatorname{ad} x \text { ad } y \operatorname{ad} z)-\operatorname{tr}(\operatorname{ad} y \text { ad } x \operatorname{ad} z) \\
& =\operatorname{tr}(\operatorname{ad} x \operatorname{ad} y \operatorname{ad} z)-\operatorname{tr}(\operatorname{ad} x \operatorname{ad} z \text { ad } y) \\
& =\operatorname{tr}(\operatorname{ad} x(\operatorname{ad} y \text { ad } z-\operatorname{ad} z \text { ad } y))=\operatorname{tr}(\operatorname{ad} x \operatorname{ad}[y z])=\langle x,[y z]\rangle
\end{aligned}
$$

Proposition 4.6 Let $I$ be an ideal of $L$ and $x, y \in I$. Then

$$
\langle x, y\rangle_{I}=\langle x, y\rangle_{L}
$$

Thus the Killing form of $L$ restricted to $I$ is the Killing form of $I$.

Proof. We choose a basis of $I$ and extend it to give a basis of $L$. With respect to this basis ad $x: L \rightarrow L$ is represented by a matrix of form

$$
\left(\begin{array}{cc}
A_{1} & A_{2} \\
O & O
\end{array}\right)
$$

since $x \in I$, and similarly ad $y: L \rightarrow L$ is represented by a matrix of form

$$
\left(\begin{array}{cc}
B_{1} & B_{2} \\
O & O
\end{array}\right)
$$

Thus ad $x$ ad $y: L \rightarrow L$ is represented by the matrix

$$
\left(\begin{array}{cc}
A_{1} B_{1} & A_{1} B_{2} \\
O & O
\end{array}\right)
$$

Hence $\operatorname{tr}_{L}(\operatorname{ad} x$ ad $y)=\operatorname{tr} A_{1} B_{1}=\operatorname{tr}_{I}(\operatorname{ad} x$ ad $y)$ and so $\langle x, y\rangle_{L}=\langle x, y\rangle_{I}$
For any subspace $M$ of $L$ we define $M^{\perp}$ by

$$
M^{\perp}=\{x \in L ;\langle x, y\rangle=0 \quad \text { for all } y \in M\} .
$$

$M^{\perp}$ is also a subspace of $L$.

Lemma 4.7 If $I$ is an ideal of $L$ then $I^{\perp}$ is also an ideal of $L$.

Proof. Let $x \in I^{\perp}, y \in L$. We must show that $[x y] \in I^{\perp}$. So let $z \in I$. Then

$$
\langle[x y], z\rangle=\langle x,[y z]\rangle=0
$$

since $[y z] \in I$ and $x \in I^{\perp}$. Thus $[x y] \in I^{\perp}$ and $I^{\perp}$ is an ideal of $L$.
We see in particular that $L^{\perp}$ is an ideal of $L$. The Killing form of $L$ is said to be non-degenerate if $L^{\perp}=O$. This is equivalent to the condition that if $\langle x, y\rangle=0$ for all $y \in L$ then $x=0$.

The Killing form of $L$ is identically zero if $L^{\perp}=L$. This means that $\langle x, y\rangle=0$ for all $x, y \in L$.

We now prove a deeper result on the Killing form which will be very useful subsequently.

Proposition 4.8 Let $L$ be a Lie algebra such that $L \neq 0$ and $L^{2}=L$. Let $H$ be a Cartan subalgebra of $L$. Then there exists $x \in H$ such that $\langle x, x\rangle \neq 0$.

Proof. We consider the Cartan decomposition of $L$ with respect to $H$. Let this be $L=\oplus L_{\lambda}$. Then we have

$$
L^{2}=[L L]=\left[\bigoplus_{\lambda} L_{\lambda}, \bigoplus_{\mu} L_{\mu}\right]=\sum_{\lambda, \mu}\left[L_{\lambda} L_{\mu}\right]
$$

Now we have $\left[L_{\lambda} L_{\mu}\right] \subset L_{\lambda+\mu}$ by Proposition 4.2. Now $L_{\lambda+\mu}=O$ if $\lambda+\mu$ is not a weight. Thus each non-zero product $\left[L_{\lambda} L_{\mu}\right]$ lies in some weight space $L_{\nu}$. We consider the zero weight space $L_{0}$. Since $L^{2}=L$ we have

$$
L_{0}=\sum_{\lambda}\left[L_{\lambda} L_{-\lambda}\right]
$$

summed over all weights $\lambda$ such that $-\lambda$ is also a weight. Now $L_{0}=H$ by Proposition 4.1, thus we have

$$
H=[H H]+\sum_{\alpha}\left[L_{\alpha} L_{-\alpha}\right]
$$

summed over all roots $\alpha \in \Phi$ such that $-\alpha$ is also a root.
Now $L$ is not nilpotent since $L^{2}=L . H$ is nilpotent and so $H \neq L$. So there is at least one root $\beta \in \Phi . \beta$ is a 1 -dimensional representation of $H$ and so vanishes on $[H H]$ since

$$
\beta[x y]=\beta(x) \beta(y)-\beta(y) \beta(x)=0 \quad x, y \in H
$$

But $\beta$ does not vanish on $H$ since $\beta \neq 0$. So using the above decomposition of $H$ we see that there is some root $\alpha \in \Phi$ such that $-\alpha \in \Phi$ and $\beta$ does not vanish on $\left[L_{\alpha} L_{-\alpha}\right]$.

We choose $x \in\left[L_{\alpha} L_{-\alpha}\right]$ such that $\beta(x) \neq 0$. Then we have

$$
\langle x, x\rangle=\operatorname{tr}(\operatorname{ad} x \operatorname{ad} x)=\sum_{\lambda} \operatorname{dim} L_{\lambda}(\lambda(x))^{2}
$$

since ad $x$ is represented on $L_{\lambda}$ by a matrix of form

$$
\left(\begin{array}{cccc}
\lambda(x) & & & * \\
& \cdot & & \\
& \cdot & \\
& & \cdot & \\
O & & & \lambda(x)
\end{array}\right)
$$

Now by Proposition 4.4 there exists $r_{\lambda, \alpha} \in \mathbb{Q}$ such that $\lambda(x)=r_{\lambda, \alpha} \alpha(x)$. Thus we have

$$
\langle x, x\rangle=\left(\sum_{\lambda} \operatorname{dim} L_{\lambda} r_{\lambda, \alpha}^{2}\right) \alpha(x)^{2}
$$

Now $\beta(x)=r_{\beta, \alpha} \alpha(x)$ and $\beta(x) \neq 0$. Thus $\alpha(x) \neq 0$ and $r_{\beta, \alpha} \neq 0$. It follows that $\langle x, x\rangle \neq 0$.

We shall now obtain some important consequences of this result.

Theorem 4.9 If the Killing form of $L$ is identically zero then $L$ is soluble.

Proof. We use induction on the dimension of $L$. If $\operatorname{dim} L=1$ then $L$ is soluble. So suppose $\operatorname{dim} L>1$. By Proposition 4.8 we have $L \neq L^{2} . L^{2}$ is an ideal of $L$ so the Killing form of $L^{2}$ is the restriction of the Killing form of $L$, by Proposition 4.6. Thus the Killing form of $L^{2}$ is identically zero. By induction $L^{2}$ is soluble. Since $L / L^{2}$ is soluble it follows that $L$ is soluble, by Proposition 1.12.

Theorem 4.10 The Killing form of $L$ is non-degenerate if and only if $L$ is semisimple.

Proof. Suppose first that the Killing form of $L$ is degenerate. Then $L^{\perp} \neq O$. Now $L^{\perp}$ is an ideal of $L$ by Lemma 4.7. Thus the Killing form of $L^{\perp}$ is the restriction of that of $L$ by Proposition 4.6. Thus the Killing form of $L^{\perp}$ is identically zero. This implies that $L^{\perp}$ is soluble, by Theorem 4.9. Thus $L$ has a non-zero soluble ideal, so $L$ is not semisimple.

Now suppose conversely that $L$ is not semisimple. Then the soluble radical $R$ of $L$ is non-zero. We consider the chain

$$
R \supset R^{(1)} \supset R^{(2)} \supset \cdots \supset R^{(k-1)} \supset R^{(k)}=O
$$

where as usual $R^{(i+1)}=\left[R^{(i)} R^{(i)}\right]$. The subspaces $R^{(i)}$ are all ideals of $L$ since the product of two ideals is an ideal. Let $I=R^{(k-1)}$. Then $I$ is a non-zero ideal of $L$ such that $I^{2}=O$.

We choose a basis of $I$ and extend it to a basis of $L$. Let $x \in I$ and $y \in L$. With respect to this basis ad $x$ is represented by a matrix of form

$$
\left(\begin{array}{ll}
O & A \\
O & O
\end{array}\right)
$$

since $I^{2}=O$ and $I$ is an ideal of $L$, ad $y$ is represented by a matrix of form

$$
\left(\begin{array}{ll}
B_{1} & B_{2} \\
O & B_{3}
\end{array}\right)
$$

and ad $x$ ad $y$ is represented by the matrix

$$
\left(\begin{array}{cc}
O & A B_{3} \\
O & O
\end{array}\right)
$$

Hence $\langle x, y\rangle=\operatorname{tr}(\operatorname{ad} x$ ad $y)=0$. Since this holds for all $x \in I$ and $y \in L$ we have $I \subset L^{\perp}$. Thus $L^{\perp} \neq O$ and so the Killing form of $L$ is degenerate.

We now define the direct sum of Lie algebras $L_{1}, L_{2} . L_{1} \oplus L_{2}$ is the vector space of all pairs $\left(x_{1}, x_{2}\right)$ with $x_{1} \in L_{1}, x_{2} \in L_{2}$ under the Lie multiplication given by

$$
\left[\left(x_{1}, x_{2}\right)\left(y_{1}, y_{2}\right)\right]=\left(\left[x_{1} y_{1}\right],\left[x_{2} y_{2}\right]\right)
$$

In this direct sum we define $I_{1}=\left\{\left(x_{1}, 0\right) ; x_{1} \in L_{1}\right\}$ and $I_{2}=\left\{\left(0, x_{2}\right) ; x_{2} \in L_{2}\right\}$. Then $I_{1}$ and $I_{2}$ are ideals of $L_{1} \oplus L_{2}$ such that $I_{1} \cap I_{2}=O$ and $I_{1}+I_{2}=L_{1} \oplus L_{2}$. Moreover $I_{1}$ is isomorphic to $L_{1}$ and $I_{2}$ is isomorphic to $L_{2}$.

Conversely let $L$ be a Lie algebra containing two ideals $I_{1}, I_{2}$ such that $I_{1} \cap I_{2}=O$ and $I_{1}+I_{2}=L$. Then the Lie algebra $I_{1} \oplus I_{2}$ is isomorphic to $L$ under the isomorphism

$$
\begin{aligned}
& \theta: I_{1} \oplus I_{2} \rightarrow L \\
& \left(x_{1}, x_{2}\right) \rightarrow x_{1}+x_{2}
\end{aligned}
$$

For $\theta$ is certainly an isomorphism of vector spaces. But $\theta$ also preserves Lie multiplication. To see this we first observe that

$$
\left[I_{1} I_{2}\right] \subset I_{1} \cap I_{2}=O
$$

Thus

$$
\begin{aligned}
{\left[\theta\left(x_{1}, x_{2}\right), \theta\left(y_{1}, y_{2}\right)\right] } & =\left[x_{1}+x_{2}, y_{1}+y_{2}\right]=\left[x_{1} y_{1}\right]+\left[x_{2} y_{2}\right] \\
& =\theta\left(\left[x_{1} y_{1}\right],\left[x_{2} y_{2}\right]\right)=\theta\left[\left(x_{1}, x_{2}\right),\left(y_{1}, y_{2}\right)\right]
\end{aligned}
$$

Thus if a Lie algebra has two complementary ideals $I_{1}, I_{2}$ the Lie algebra is isomorphic to $I_{1} \oplus I_{2}$.

We may in a similar way consider direct sums $L_{1} \oplus L_{2} \oplus \cdots \oplus L_{n}$ of more than two Lie algebras.

Theorem 4.11 A Lie algebra $L$ is semisimple if and only if $L$ is isomorphic to a direct sum of non-trivial simple Lie algebras.

Proof. Suppose $L$ is semisimple. If $L$ is simple then $L$ must be non-trivial since the trivial simple Lie algebra is not semisimple. Thus we suppose $L$ is not simple. Let $I$ be a minimal non-zero ideal of $L$. Then $I \neq O$ and $I \neq L$. Consider the subspace $I^{\perp}$ of $L ; I^{\perp}$ is also an ideal of $L$ by Lemma 4.7. Now the Killing form of $L$ is non-degenerate by Theorem 4.10. Thus an element $x \in L$ lies in $I^{\perp}$ if and only if the coordinates of $x$ with respect to a basis of $L$ satisfy $\operatorname{dim} I$ homogeneous linear equations which are linearly independent. It follows that

$$
\operatorname{dim} I^{\perp}=\operatorname{dim} L-\operatorname{dim} I
$$

Now consider the subspace $I \cap I^{\perp}$. This is an ideal of $L$. Thus the Killing form of $I \cap I^{\perp}$ is the restriction of the Killing form of $L$, by Proposition 4.6. Hence $I \cap I^{\perp}$ is soluble, by Theorem 4.9. Since $L$ is semisimple we have $I \cap I^{\perp}=O$. Thus

$$
\begin{aligned}
\operatorname{dim}\left(I+I^{\perp}\right) & =\operatorname{dim} I+\operatorname{dim} I^{\perp}-\operatorname{dim}\left(I \cap I^{\perp}\right) \\
& =\operatorname{dim} I+\operatorname{dim} I^{\perp}=\operatorname{dim} L
\end{aligned}
$$

Hence $I+I^{\perp}=L$. Thus $L$ is the direct sum of its ideals $I$ and $I^{\perp}$. Hence $L$ is isomorphic to the Lie algebra $I \oplus I^{\perp}$.

We shall now show that $I$ is a simple Lie algebra. Let $J$ be an ideal of $I$. Then we have

$$
[J L] \subset[J I]+\left[J I^{\perp}\right] \subset[J I] \subset J
$$

since $\left[J I^{\perp}\right] \subset\left[I I^{\perp}\right] \subset I \cap I^{\perp}=O$. Thus $J$ is an ideal of $L$ contained in $I$. Since $I$ is a minimal ideal of $L$ we have $J=O$ or $J=I$. Thus $I$ is simple.

We show next that $I^{\perp}$ is semisimple. Let $J$ be a soluble ideal of $I^{\perp}$. Then

$$
[J L] \subset[J I]+\left[J I^{\perp}\right] \subset\left[J I^{\perp}\right] \subset J
$$

since $[J I] \subset\left[I^{\perp} I\right] \subset I \cap I^{\perp}=O$. Thus $J$ is an ideal of $L$. Since $L$ is semisimple and $J$ is soluble we have $J=O$. Thus $I^{\perp}$ is semisimple.

Now we know $\operatorname{dim} I^{\perp}<\operatorname{dim} L$. By induction we may assume $I^{\perp}$ is a direct sum of simple non-trivial Lie algebras. Since $L=I \oplus I^{\perp}$ and $I$ is simple and non-trivial, $L$ is also a direct sum of simple non-trivial Lie algebras.

Conversely suppose that

$$
L=L_{1} \oplus \cdots \oplus L_{r}
$$

where each $L_{i}$ is a simple non-trivial Lie algebra. Each $L_{i}$ is semisimple so has non-degenerate Killing form by Theorem 4.10. Now each $L_{i}$ is an ideal of $L$. Moreover if $x_{i} \in L_{i}, x_{j} \in L_{j}$ and $i \neq j$ then $\left\langle x_{i}, x_{j}\right\rangle=0$. For

$$
\text { ad } x_{i} \text { ad } x_{j} \cdot y \in L_{i} \cap L_{j}=O \quad \text { for all } y \in L
$$

thus $\left\langle x_{i}, x_{j}\right\rangle=\operatorname{tr}\left(\operatorname{ad} x_{i}\right.$ ad $\left.x_{j}\right)=0$.
Now let $x=x_{1}+\cdots+x_{r} \in L^{\perp}$ with $x_{i} \in L_{i}$. Let $y_{i} \in L_{i}$. Then we have

$$
\left\langle x_{i}, y_{i}\right\rangle=\left\langle x, y_{i}\right\rangle=0
$$

Since this holds for all $y_{i} \in L_{i}$ we have $x_{i}=0$. This holds for all $i$, hence $x=0$. Thus $L^{\perp}=O$ and the Killing form of $L$ is non-degenerate. This implies that $L$ is semisimple by Theorem 4.10.

### 4.3 The Cartan decomposition of a semisimple Lie algebra

When $L$ is semisimple we can say much more about its Cartan decomposition than in the general case. We shall now investigate this Cartan decomposition in detail.

Let $L$ be semisimple, $H$ be a Cartan subalgebra of $L$, and $L=\bigoplus L_{\lambda}$ be the Cartan decomposition of $L$ with respect to $H$. We recall from Proposition 4.1 that $L_{0}=H$.

Proposition 4.12 $L_{\lambda}$ and $L_{\mu}$ are orthogonal with respect to the Killing form, provided $\mu \neq-\lambda$.

Proof. Let $x \in L_{\lambda}, y \in L_{\mu}$. We assume $\lambda+\mu \neq 0$ and must show that $\langle x, y\rangle=0$. Now for any weight space $L_{\nu}$ we have

$$
\operatorname{ad} x \text { ad } y L_{\nu} \subset L_{\lambda+\mu+\nu} \quad \text { by Proposition 4.2 }
$$

We choose a basis of $L$ adapted to the Cartan decomposition. With respect to such a basis ad $x$ ad $y$ will be represented by a block matrix of form

$$
\left(\begin{array}{lllllll}
0 & & & & & & \\
& 0 & & & & * & \\
& & \cdot & & & & \\
& & & \cdot & & & \\
& & & & \cdot & & \\
& * & & & & 0 & \\
& & & & & & 0
\end{array}\right)
$$

since $\lambda+\mu+\nu \neq \nu$. Hence we have

$$
\langle x, y\rangle=\operatorname{tr}(\operatorname{ad} x \text { ad } y)=0
$$

Proposition 4.13 If $\alpha$ is a root of $L$ with respect to $H$ then $-\alpha$ is also a root.

Proof. We recall that $\alpha$ is a root if $\alpha \neq 0$ and $L_{\alpha} \neq O$. Suppose if possible that $-\alpha$ is not a root. Since $-\alpha \neq 0$ we have $L_{-\alpha}=O$. By Proposition 4.12 we see that $L_{\alpha}$ is orthogonal to all $L_{\lambda}$, hence $L_{\alpha} \subset L^{\perp}$. But since $L$ is semisimple we have $L^{\perp}=O$ by Theorem 4.10. Thus $L_{\alpha}=O$, which contradicts the fact that $\alpha$ is a root.

Proposition 4.14 The Killing form of L remains non-degenerate on restriction to $H$. Thus if $x \in H$ satisfies $\langle x, y\rangle=0$ for all $y \in H$ then $x=0$.

Proof. Let $x \in H$ and suppose $\langle x, y\rangle=0$ for all $y \in H$. We also have $\langle x, y\rangle=0$ for all $y \in L_{\alpha}$ where $\alpha \neq 0$, by Proposition 4.12. Thus $\langle x, y\rangle=0$ for all $y \in L$ and so $x \in L^{\perp}$. Since $L$ is semisimple $L^{\perp}=O$, hence $x=0$ as required.

Note that the Killing form of $L$ restricted to $H$ does not coincide with the Killing form of $H$. The latter is degenerate since $H$ is not semisimple.

Theorem $4.15[H H]=O$. Thus the Cartan subalgebras of a semisimple Lie algebra are abelian.

Proof. Let $x \in[H H]$ and $y \in H$. Then we have

$$
\langle x, y\rangle=\operatorname{tr}(\operatorname{ad} x \operatorname{ad} y)=\sum_{\lambda} \operatorname{dim} L_{\lambda} \lambda(x) \lambda(y)
$$

since ad $x$ ad $y$ is represented on $L_{\lambda}$ by a matrix of form

$$
\left(\begin{array}{cccc}
\lambda(x) \lambda(y) & & & * \\
& \cdot & & \\
O & \cdot & \\
& & \cdot & \\
& & & \lambda(x) \lambda(y)
\end{array}\right)
$$

However, $\lambda$ is a 1 -dimensional representation of $H$ and $x \in[H H]$, hence $\lambda(x)=0$. Thus $\langle x, y\rangle=0$ for all $y \in H$. This implies $x=0$ by Proposition 4.14. Thus $[H H]=O$.

Let $H^{*}=\operatorname{Hom}(H, \mathbb{C})$ be the dual space of $H$. This is the vector space of all linear maps from $H$ to $\mathbb{C}$. We have $\operatorname{dim} H^{*}=\operatorname{dim} H$.

We define a map $H \rightarrow H^{*}$ using the Killing form of $L$. Given $h \in H$ we define $h^{*} \in H^{*}$ by

$$
h^{*}(x)=\langle h, x\rangle \quad \text { for all } x \in H
$$

Lemma 4.16 The map $h \rightarrow h^{*}$ is an isomorphism of vector spaces between $H$ and $H^{*}$.

Proof. The map is certainly linear. Suppose $h \in H$ lies in the kernel. Then $\langle h, x\rangle=0$ for all $x \in H$. This implies $h=0$ by Proposition 4.14. Thus the kernel is $O$. Hence the image must be the whole of $H^{*}$, since $\operatorname{dim} H^{*}=\operatorname{dim} H$. Hence our map is bijective.

Now we have a finite subset $\Phi \subset H^{*}$, the set of roots of $L$ with respect to $H$. For each $\alpha \in \Phi$ there is a unique element $h_{\alpha}^{\prime} \in H$ such that

$$
\alpha(x)=\left\langle h_{\alpha}^{\prime}, x\right\rangle \quad \text { for all } x \in H
$$

(The notation $h_{\alpha}$ might seem more natural, but this will be reserved for the coroot of $\alpha$, to be discussed in Chapter 7.)

Proposition 4.17 The vectors $h_{\alpha}^{\prime}$ for $\alpha \in \Phi$ span $H$.

Proof. Suppose if possible that the $h_{\alpha}^{\prime}$ lie in a proper subspace of $H$. Then there exists an element $x \in H$ with $x \neq 0$ and $\left\langle h_{\alpha}^{\prime}, x\right\rangle=0$ for all $\alpha \in \Phi$. Thus $\alpha(x)=0$ for all $\alpha \in \Phi$. Let $y \in H$. Then we have

$$
\langle x, y\rangle=\operatorname{tr}(\operatorname{ad} x \text { ad } y)=\sum_{\lambda} \operatorname{dim} L_{\lambda} \lambda(x) \lambda(y)=0
$$

since $\lambda(x)=0$ for all weights $\lambda$. Thus $\langle x, y\rangle=0$ for all $y \in H$. This implies $x=0$ by Proposition 4.14, a contradiction.

Proposition $4.18 h_{\alpha}^{\prime} \in\left[L_{\alpha} L_{-\alpha}\right]$ for all $\alpha \in \Phi$.

Proof. $L_{\alpha}$ is an $H$-module. Since all irreducible $H$-modules are 1-dimensional $L_{\alpha}$ contains a 1-dimensional $H$-submodule $\mathbb{C} e_{\alpha}$. We have $\left[x e_{\alpha}\right]=\alpha(x) e_{\alpha}$ for all $x \in H$.

Let $y \in L_{-\alpha}$. Then $\left[e_{\alpha} y\right] \in\left[L_{\alpha} L_{-\alpha}\right] \subset H$. We shall show that $\left[e_{\alpha} y\right]=$ $\left\langle e_{\alpha}, y\right\rangle h_{\alpha}^{\prime}$. In order to prove this we define

$$
z=\left[e_{\alpha} y\right]-\left\langle e_{\alpha}, y\right\rangle h_{\alpha}^{\prime} \in H
$$

Let $x \in H$. Then

$$
\begin{aligned}
\langle x, z\rangle & =\left\langle x,\left[e_{\alpha} y\right]\right\rangle-\left\langle e_{\alpha}, y\right\rangle\left\langle x, h_{\alpha}^{\prime}\right\rangle \\
& =\left\langle\left[x e_{\alpha}\right], y\right\rangle-\left\langle e_{\alpha}, y\right\rangle \alpha(x) \\
& =\alpha(x)\left\langle e_{\alpha}, y\right\rangle-\left\langle e_{\alpha}, y\right\rangle \alpha(x)=0 .
\end{aligned}
$$

Thus $\langle x, z\rangle=0$ for all $x \in H$, and it follows that $z=0$. Hence

$$
\left[e_{\alpha} y\right]=\left\langle e_{\alpha}, y\right\rangle h_{\alpha}^{\prime} \quad \text { for all } \quad y \in L_{-\alpha} .
$$

Now we can choose $y \in L_{-\alpha}$ such that $\left\langle e_{\alpha}, y\right\rangle \neq 0$. Otherwise $e_{\alpha}$ would be orthogonal to $L_{-\alpha}$, so orthogonal to the whole of $L$ by Proposition 4.12. Then $e_{\alpha} \in L^{\perp}$. But $L^{\perp}=0$ since $L$ is semisimple. Thus $e_{\alpha}=0$, a contradiction. Thus we can find $y \in L_{-\alpha}$ with $\left\langle e_{\alpha}, y\right\rangle \neq 0$. Then

$$
h_{\alpha}^{\prime}=\frac{1}{\left\langle e_{\alpha}, y\right\rangle}\left[e_{\alpha} y\right] \in\left[L_{\alpha} L_{-\alpha}\right] .
$$

Proposition $4.19\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle \neq 0$ for all $\alpha \in \Phi$.

Proof. We suppose that $\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle=0$ for some $\alpha \in \Phi$ and obtain a contradiction. Let $\beta$ be any element of $\Phi$. By Proposition 4.4 there is a number $r_{\beta, \alpha} \in \mathbb{Q}$ such that $\beta=r_{\beta, \alpha} \alpha$ when restricted to $\left[L_{\alpha} L_{-\alpha}\right]$. Since $h_{\alpha}^{\prime} \in\left[L_{\alpha} L_{-\alpha}\right]$ by Proposition 4.18 we obtain

$$
\beta\left(h_{\alpha}^{\prime}\right)=r_{\beta, \alpha} \alpha\left(h_{\alpha}^{\prime}\right)
$$

that is $\left\langle h_{\beta}^{\prime}, h_{\alpha}^{\prime}\right\rangle=r_{\beta, \alpha}\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle=0$.
This holds for all $\beta \in \Phi$. But by Proposition 4.17 the elements $h_{\beta}^{\prime}$ for $\beta \in \Phi$ span $H$. Thus we have $\left\langle x, h_{\alpha}^{\prime}\right\rangle=0$ for all $x \in H$. This implies that $h_{\alpha}^{\prime}=0$ by Proposition 4.14. This in turn implies that $\alpha=0$, which contradicts $\alpha \in \Phi$.

Having obtained a number of results on the Cartan decomposition of a semisimple Lie algebra, each depending on previous results, we are now able to obtain one of the most important properties of the Cartan decomposition.

Theorem 4.20 $\operatorname{dim} L_{\alpha}=1$ for all $\alpha \in \Phi$.

Proof. We choose a 1-dimensional $H$-submodule $\mathbb{C} e_{\alpha}$ of $L_{\alpha}$ as in Proposition 4.18 and, as in the proof of that proposition, we can find an element $e_{-\alpha} \in L_{-\alpha}$ with $\left[e_{\alpha} e_{-\alpha}\right]=h_{\alpha}^{\prime}$.

We consider the subspace $M$ of $L$ given by

$$
M=\mathbb{C} e_{\alpha} \oplus \mathbb{C} h_{\alpha}^{\prime} \oplus L_{-\alpha} \oplus L_{-2 \alpha} \oplus \cdots
$$

There are only finitely many summands of $M$ since $\Phi$ is finite and there are only finitely many non-negative integers $r$ with $L_{-r \alpha} \neq O$.

We observe that ad $e_{\alpha} M \subset M$. For

$$
\begin{aligned}
& {\left[e_{\alpha} e_{\alpha}\right]=0} \\
& {\left[e_{\alpha} h_{\alpha}^{\prime}\right]=-\alpha\left(h_{\alpha}^{\prime}\right) e_{\alpha}} \\
& {\left[e_{\alpha} y\right]=\left\langle e_{\alpha}, y\right\rangle h_{\alpha}^{\prime} \quad \text { for all } y \in L_{-\alpha}}
\end{aligned}
$$

by the proof of Proposition 4.18, and

$$
\text { ad } e_{\alpha} \cdot L_{-r \alpha} \subset L_{-(r-1) \alpha} \quad \text { for all } r \geq 2
$$

by Proposition 4.2.

Similarly we can show that ad $e_{-\alpha} M \subset M$. For we have

$$
\begin{aligned}
& {\left[e_{-\alpha} e_{\alpha}\right]=-h_{\alpha}^{\prime}} \\
& {\left[e_{-\alpha} h_{\alpha}^{\prime}\right]=\alpha\left(h_{\alpha}^{\prime}\right) e_{-\alpha}}
\end{aligned}
$$

and ad $e_{-\alpha} L_{-r \alpha} \subset L_{-(r+1) \alpha}$ for all $r \geq 1$.
Now $h_{\alpha}^{\prime}=\left[e_{\alpha} e_{-\alpha}\right]$ and so

$$
\operatorname{ad} h_{\alpha}^{\prime}=\operatorname{ad} e_{\alpha} \operatorname{ad} e_{-\alpha}-\operatorname{ad} e_{-\alpha} \operatorname{ad} e_{\alpha}
$$

Hence ad $h_{\alpha}^{\prime} M \subset M$. We shall calculate the trace of ad $h_{\alpha}^{\prime}$ on $M$ in two different ways. On the one hand we have

$$
\begin{aligned}
\operatorname{tr}_{M}\left(\operatorname{ad} h_{\alpha}^{\prime}\right) & =\alpha\left(h_{\alpha}^{\prime}\right)+\operatorname{dim} L_{-\alpha}\left(-\alpha\left(h_{\alpha}^{\prime}\right)\right)+\operatorname{dim} L_{-2 \alpha}\left(-2 \alpha\left(h_{\alpha}^{\prime}\right)\right)+\cdots \\
& =\alpha\left(h_{\alpha}^{\prime}\right)\left(1-\operatorname{dim} L_{-\alpha}-2 \operatorname{dim} L_{-2 \alpha}-\cdots\right) .
\end{aligned}
$$

On the other hand we have

$$
\operatorname{tr}_{M}\left(\operatorname{ad} h_{\alpha}^{\prime}\right)=\operatorname{tr}_{M}\left(\operatorname{ad} e_{\alpha} \text { ad } e_{-\alpha}-\operatorname{ad} e_{-\alpha} \text { ad } e_{\alpha}\right)=0
$$

Thus

$$
\alpha\left(h_{\alpha}^{\prime}\right)\left(1-\operatorname{dim} L_{-\alpha}-2 \operatorname{dim} L_{-2 \alpha}-\cdots\right)=0 .
$$

Now $\alpha\left(h_{\alpha}^{\prime}\right)=\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle \neq 0$ by Proposition 4.19. Thus

$$
1-\operatorname{dim} L_{-\alpha}-2 \operatorname{dim} L_{-2 \alpha}-\cdots=0
$$

This implies that $\operatorname{dim} L_{-\alpha}=1$ and $\operatorname{dim} L_{-r \alpha}=0$ for all $r \geq 2$. Now $\alpha \in \Phi$ if and only if $-\alpha \in \Phi$, by Proposition 4.13. Thus $\operatorname{dim} L_{\alpha}=1$ for all $\alpha \in \Phi$.

Note that although all the root spaces $L_{\alpha}$ are 1-dimensional the space $H=L_{0}$ need not be 1-dimensional.

Proposition 4.21 If $\alpha \in \Phi$ and $r \alpha \in \Phi$ where $r \in \mathbb{Z}$ then $r=1$ or -1 .

Proof. This follows from the proof of Theorem 4.20, where we showed that, for all $\alpha \in \Phi,-r \alpha \notin \Phi$ for all $r \geq 2$. This, together with the fact that $r \alpha \in \Phi$ if and only if $-r \alpha \in \Phi$, gives the required result.

We shall now obtain some further properties of the set $\Phi$ of roots. Let $\alpha, \beta \in \Phi$ be such that $\beta \neq \alpha$ and $\beta \neq-\alpha$. Then $\beta$ cannot be an integer multiple of $\alpha$, by Proposition 4.21. There exist integers $p \geq 0, q \geq 0$ such that the elements

$$
-p \alpha+\beta, \ldots,-\alpha+\beta, \beta, \alpha+\beta, \ldots, q \alpha+\beta
$$

all lie in $\Phi$, but $-(p+1) \alpha+\beta$ and $(q+1) \alpha+\beta$ do not lie in $\Phi$. The set of roots

$$
-p \alpha+\beta, \ldots, q \alpha+\beta
$$

is called the $\boldsymbol{\alpha}$-chain of roots through $\beta$. Let $M$ be the subspace of $L$ defined by

$$
M=L_{-p \alpha+\beta} \oplus \cdots \oplus L_{q \alpha+\beta}
$$

Then we have ad $e_{\alpha} M \subset M$. This follows from the fact that ad $e_{\alpha} L_{r \alpha+\beta} \subset$ $L_{(r+1) \alpha+\beta}$ and $L_{(q+1) \alpha+\beta}=0$ since $(q+1) \alpha+\beta \notin \Phi$ and $(q+1) \alpha+\beta \neq 0$. Similarly we see that ad $e_{-\alpha} M \subset M$.

We assume that $\left[e_{\alpha} e_{-\alpha}\right]=h_{\alpha}^{\prime}$, as in the proof of Theorem 4.20. Then we have

$$
\operatorname{ad} h_{\alpha}^{\prime}=\operatorname{ad} e_{\alpha} \operatorname{ad} e_{-\alpha}-\operatorname{ad} e_{-\alpha} \operatorname{ad} e_{\alpha}
$$

and so ad $h_{\alpha}^{\prime}, M \subset M$. We calculate the trace of ad $h_{\alpha}^{\prime}$ on $M$ in two different ways. We have

$$
\operatorname{tr}_{M}\left(\operatorname{ad} h_{\alpha}^{\prime}\right)=\sum_{r=-p}^{q}(r \alpha+\beta)\left(h_{\alpha}^{\prime}\right)
$$

since $\operatorname{dim} L_{r \alpha+\beta}=1$. We also have

$$
\operatorname{tr}_{M}\left(\operatorname{ad} h_{\alpha}^{\prime}\right)=\operatorname{tr}_{M}\left(\operatorname{ad} e_{\alpha} \operatorname{ad} e_{-\alpha}\right)-\operatorname{tr}_{M}\left(\operatorname{ad} e_{-\alpha} \operatorname{ad} e_{\alpha}\right)=0
$$

Thus

$$
\sum_{r=-p}^{q}(r \alpha+\beta)\left(h_{\alpha}^{\prime}\right)=0
$$

that is

$$
\left(\frac{q(q+1)}{2}-\frac{p(p+1)}{2}\right) \alpha\left(h_{\alpha}^{\prime}\right)+(p+q+1) \beta\left(h_{\alpha}^{\prime}\right)=0 .
$$

Since $p+q+1 \neq 0$ we obtain

$$
\frac{(q-p)}{2}\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle+\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle=0
$$

that is

$$
2 \frac{\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}=p-q
$$

since $\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle \neq 0$ by Proposition 4.19. Thus we have proved the following result.

Proposition 4.22 Let $\alpha, \beta$ be roots such that $\beta \neq \alpha$ and $\beta \neq-\alpha$. Let

$$
-p \alpha+\beta, \ldots, \beta, \ldots, q \alpha+\beta
$$

be the $\alpha$-chain of roots through $\beta$. Then

$$
2 \frac{\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}=p-q
$$

This result has some useful corollaries. The first gives a strengthening of the result of Proposition 4.21.

Proposition 4.23 If $\alpha \in \Phi$ and $\zeta \alpha \in \Phi$ where $\zeta \in \mathbb{C}$, then $\zeta=1$ or -1 .

Proof. Suppose if possible that $\zeta \neq \pm 1$. We put $\beta=\zeta \alpha$ and apply Proposition 4.22. This gives

$$
2 \zeta=2 \frac{\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}=p-q
$$

Hence $2 \zeta \in \mathbb{Z}$. If $\zeta \in \mathbb{Z}$ then $\zeta= \pm 1$ by Proposition 4.21. Hence $\zeta \notin \mathbb{Z}$. Then the $\alpha$-chain of roots through $\beta$ is

$$
-\left(\frac{p+q}{2}\right) \alpha, \ldots, \beta=\left(\frac{p-q}{2}\right) \alpha, \ldots,\left(\frac{p+q}{2}\right) \alpha
$$

Now $p, q$ are not both 0 since $\beta \neq 0$. So all the roots in the $\alpha$-chain are odd multiples of $\frac{1}{2} \alpha$. Since the first and the last are negatives of one another and consecutive roots differ by $\alpha$ it is clear that $\frac{1}{2} \alpha$ lies in the chain. Hence $\frac{1}{2} \alpha \in \Phi$. Since $\alpha \in \Phi$ we have a contradiction to Proposition 4.21. Hence $\zeta$ must be 1 or -1 .

Thus the only roots which are scalar multiples of a root $\alpha$ are $\alpha$ and $-\alpha$.
Proposition $4.24\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle \in \mathbb{Q}$ for all $\alpha, \beta \in \Phi$.
Proof. We know from the outset that $\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle \in \mathbb{C}$. Now we have

$$
2 \frac{\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle} \in \mathbb{Z} \quad \text { by Proposition } 4.22
$$

Thus $\frac{\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle} \in \mathbb{Q}$. It will therefore be sufficient to show that $\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle \in \mathbb{Q}$. Now we have

$$
\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle=\operatorname{tr}\left(\operatorname{ad} h_{\alpha}^{\prime} \operatorname{ad} h_{\alpha}^{\prime}\right)=\sum_{\beta \in \Phi}\left(\beta\left(h_{\alpha}^{\prime}\right)\right)^{2}=\sum_{\beta \in \Phi}\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle^{2} .
$$

If follows that

$$
\frac{1}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}=\sum_{\beta \in \Phi}\left(\frac{\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}\right)^{2} \in \mathbb{Q}
$$

Hence $\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle \in \mathbb{Q}$ and the result is proved.

### 4.4 The Lie algebra $\mathfrak{E l}_{n}(\mathbb{C})$

We shall now illustrate the general results about the Cartan decomposition of a semisimple Lie algebra by considering in detail the Lie algebra $\mathfrak{E l}_{n}(\mathbb{C})$. The special linear Lie algebra $\mathfrak{E l}_{n}(\mathbb{C})$ is the Lie algebra of all $n \times n$ matrices of trace 0 under Lie multiplication $[A B]=A B-B A . \mathfrak{F l}_{n}(\mathbb{C})$ is a subalgebra of $\mathfrak{g l}_{n}(\mathbb{C})=\left[M_{n}(\mathbb{C})\right]$. We have

$$
\operatorname{dim} \mathfrak{g l}_{n}(\mathbb{C})=n^{2}, \quad \operatorname{dim} \mathfrak{F l}_{n}(\mathbb{C})=n^{2}-1
$$

We shall assume $n \geq 2$. Then $\mathfrak{s l}(\mathbb{C})$ has a basis

$$
\mathrm{E}_{11}-\mathrm{E}_{22}, \quad \mathrm{E}_{22}-\mathrm{E}_{33}, \quad \ldots, \quad \mathrm{E}_{n-1, n-1}-\mathrm{E}_{n n}, \quad \mathrm{E}_{i j} \quad i \neq j
$$

where the $\mathrm{E}_{i j}$ are elementary matrices.
Theorem $4.25 \mathfrak{g l}_{n}(\mathbb{C})$ is a simple Lie algebra.
Proof. We have $\mathfrak{g l}_{n}(\mathbb{C})=\mathfrak{E l}_{n}(\mathbb{C}) \oplus \mathbb{C} I_{n}$. Now every ideal of $\mathfrak{s l} l_{n}(\mathbb{C})$ is an ideal of $\mathfrak{g l} l_{n}(\mathbb{C})$. For $\left[I, \mathfrak{g l}_{n}(\mathbb{C})\right] \subset I$ implies $\left[I, \mathfrak{g l}_{n}(\mathbb{C})\right] \subset I$ since $\left[x, I_{n}\right]=0$ for all $x \in I$. It will therefore be sufficient to show that the only non-zero ideal of $\mathfrak{g l}_{n}(\mathbb{C})$ contained in $\mathfrak{s l}(\mathbb{C})$ is equal to $\mathfrak{S l}_{n}(\mathbb{C})$.

Let $I$ be a non-zero ideal of $\mathfrak{g l}_{n}(\mathbb{C})$ contained in $\mathfrak{F l}(\mathbb{C})$. Let $x \in I$ with $x \neq 0$. Then

$$
x=\sum x_{p q} E_{p q} \quad \text { with } x_{p q} \in \mathbb{C}
$$

Not all $x_{p q}$ are zero.
Suppose first that there exist $i \neq j$ with $x_{i j} \neq 0$. Then

$$
\left[E_{i i}, \sum x_{p q} E_{p q}\right]=\sum_{q} x_{i q} E_{i q}-\sum_{p} x_{p i} E_{p i} \in I
$$

Also

$$
\left[\left[E_{i i}, x\right], E_{j j}\right]=x_{i j} E_{i j}+x_{j i} E_{j i} \in I
$$

Hence

$$
\left[E_{i i}-E_{j j}, x_{i j} E_{i j}+x_{j i} E_{j i}\right]=2 x_{i j} E_{i j}-2 x_{j i} E_{j i} \in I
$$

Thus $4 x_{i j} E_{i j} \in I$. Since $x_{i j} \neq 0$ we have $E_{i j} \in I$.
Now suppose that $x_{i j}=0$ for all $i \neq j$. Then $x=\sum x_{p p} E_{p p}$. Since $\sum x_{p p}=0$ and not all $x_{p p}=0$ the $x_{p p}$ are not all equal. Suppose $x_{i i} \neq x_{j j}$. Then

$$
\left[x, E_{i j}\right]=\left(x_{i i}-x_{j j}\right) E_{i j} \in I
$$

and so $E_{i j} \in I$.
Thus in either case there exist $i \neq j$ with $E_{i j} \in I$. Let $q \neq i, j$. Then

$$
\left[E_{i j}, E_{j q}\right]=E_{i q} \in I
$$

Thus $E_{i q} \in I$ for all $q \neq i$. Now let $p \neq i, q$. Then

$$
\left[E_{p i}, E_{i q}\right]=E_{p q} \in I
$$

Hence $E_{p q} \in I$ for all $p \neq q$. Also

$$
\left[E_{p q}, E_{q p}\right]=E_{p p}-E_{q q} \in I \quad \text { for all } p \neq q
$$

But the $E_{p p}-E_{q q}$ for $p \neq q$ and the $E_{p q}$ for $p \neq q$ generate $\mathfrak{s l}(\mathbb{C})$. Thus $I=\mathfrak{S l}_{n}(\mathbb{C})$ and $\mathfrak{S l}_{n}(\mathbb{C})$ is simple.

We next determine a Cartan subalgebra of $\mathfrak{g l}_{n}(\mathbb{C})$. We write $L=\mathfrak{\xi l}_{n}(\mathbb{C})$.

Proposition 4.26 Let $H$ be the set of diagonal matrices in $L$. Then $\operatorname{dim} H=$ $n-1$ and $H$ is a Cartan subalgebra of $L$.

Proof. The vector space of diagonal $n \times n$ matrices of trace 0 clearly has dimension $n-1$. It is a subalgebra $H$ of $L$ with $[H H]=O$. Thus $H$ is nilpotent. To show $H$ is a Cartan subalgebra we must show $H=N(H)$.

Let $\sum_{i, j} \lambda_{i j} E_{i j}$ lie in $N(H)$. Suppose if possible that $\lambda_{i j} \neq 0$ for some $i \neq j$. We have

$$
\left[\sum_{k} \mu_{k} E_{k k}, \sum_{i, j} \lambda_{i j} E_{i j}\right] \in H
$$

for all $\sum_{k} \mu_{k} E_{k k} \in H$. The coefficient of $E_{i j}$ in this matrix is $\left(\mu_{i}-\mu_{j}\right) \lambda_{i j}$. Thus if we choose $(i, j)$ such that $i \neq j$ and $\lambda_{i j} \neq 0$ and choose $\sum_{k} \mu_{k} E_{k k} \in H$ with $\mu_{i} \neq \mu_{j}$ we obtain a contradiction. Hence $\lambda_{i j}=0$ for all $i \neq j$. Thus $N(H)=H$ and $H$ is a Cartan subalgebra of $L$.

We next obtain the Cartan decomposition of $L$ with respect to $H$.

Proposition 4.27 Let $H$ be the subalgebra of diagonal matrices in $L$. Then the Cartan decomposition of $L$ with respect to $H$ is

$$
L=H \oplus \sum_{i \neq j} \mathbb{C} E_{i j}
$$

Proof. This is certainly a decomposition of $L$ into a direct sum of subspaces. To show it is a Cartan decomposition it is sufficient to verify that the 1 -dimensional subspaces $\mathbb{C} E_{i j}$ for $i \neq j$ are $H$-submodules of $L$. Now we have

$$
\left[\sum_{k=1}^{n} \lambda_{k} E_{k k}, E_{i j}\right]=\left(\lambda_{i}-\lambda_{j}\right) E_{i j}
$$

and so $\mathbb{C} E_{i j}$ is indeed an $H$-submodule.
We next obtain the roots of $L$ with respect to $H$.

Proposition 4.28 The roots of $L$ with respect to $H$ are the functions $H \rightarrow \mathbb{C}$ given by

$$
\left(\begin{array}{cccc}
\lambda_{1} & & & O \\
& \cdot & & \\
& \cdot & & \\
& & \cdot & \\
O & & & \\
& & & \\
& & & \lambda_{n}
\end{array}\right) \rightarrow \lambda_{i}-\lambda_{j} \quad i \neq j
$$

Proof. This follows from the Cartan decomposition given in Proposition 4.27.

We next calculate the value of the Killing form $\langle x, y\rangle$ when $x, y \in H$.

Proposition 4.29 Let $x=\sum_{i=1}^{n} \lambda_{i} E_{i i}, y=\sum_{i=1}^{n} \mu_{i} E_{i i}$ lie in H. Then $\langle x, y\rangle=$ $2 n \operatorname{tr}(x y)$.

Proof. We have

$$
\langle x, y\rangle=\operatorname{tr}(\operatorname{ad} x \operatorname{ad} y)=\sum_{\substack{i, j \\ i \neq j}}\left(\lambda_{i}-\lambda_{j}\right)\left(\mu_{i}-\mu_{j}\right)
$$

since ad $x$ ad $y E_{i j}=\left(\lambda_{i}-\lambda_{j}\right)\left(\mu_{i}-\mu_{j}\right) E_{i j}$ for $i \neq j$, and ad $x$ ad $y H=O$.

Hence

$$
\begin{aligned}
\langle x, y\rangle & =\sum_{i, j}\left(\lambda_{i}-\lambda_{j}\right)\left(\mu_{i}-\mu_{j}\right) \\
& =\sum_{i, j} \lambda_{i} \mu_{i}+\sum_{i, j} \lambda_{j} \mu_{j}-\sum_{i, j} \lambda_{i} \mu_{j}-\sum_{i, j} \lambda_{j} \mu_{i} \\
& =2 n \operatorname{tr}(x y)-\left(\sum_{i} \lambda_{i}\right)\left(\sum_{j} \mu_{j}\right)-\left(\sum_{j} \lambda_{j}\right)\left(\sum_{i} \mu_{i}\right) \\
& =2 n \operatorname{tr}(x y), \quad \text { since } \sum_{i} \lambda_{i}=\sum_{i} \mu_{i}=0
\end{aligned}
$$

We may use this knowledge of the Killing form of $L$ restricted to $H$ to determine the elements $h_{\alpha}^{\prime} \in H$ corresponding to the roots $\alpha \in \Phi$.

Proposition 4.30 Let $\alpha_{i j} \in \Phi$ satisfy

$$
\alpha_{i j}\left(\begin{array}{cccc}
\lambda_{1} & & & \\
& \cdot & & O \\
& \cdot & \\
O & & \cdot & \\
& & & \lambda_{n}
\end{array}\right)=\lambda_{i}-\lambda_{j} \quad i \neq j
$$

Then $h_{\alpha_{i j}}^{\prime}=\frac{1}{2 n}\left(E_{i i}-E_{j j}\right)$.
Proof. Let $x=\sum_{k=1}^{n} \lambda_{k} E_{k k} \in H$. Then we have

$$
\begin{aligned}
\left\langle\frac{1}{2 n}\left(E_{i i}-E_{j j}\right), x\right\rangle & =2 n \operatorname{tr}\left(\frac{1}{2 n}\left(E_{i i}-E_{j j}\right) x\right) \\
& =\lambda_{i}-\lambda_{j}=\alpha_{i j}(x), \quad \text { by Proposition 4.29. }
\end{aligned}
$$

However, $h_{\alpha_{i j}}^{\prime} \in H$ is uniquely determined by the condition $\left\langle h_{\alpha_{i j}}^{\prime}, x\right\rangle=\alpha_{i j}(x)$ for all $x \in H$. Hence $h_{\alpha_{i j}}^{\prime}=\frac{1}{2 n}\left(E_{i i}-E_{j j}\right)$.

## 5

## The root system and the Weyl group

### 5.1 Positive systems and fundamental systems of roots

As before, let $L$ be a semisimple Lie algebra and $H$ be a Cartan subalgebra. Let $\Phi$ be the set of roots of $L$ with respect to $H$. We know by Proposition 4.17 that the elements $h_{\alpha}^{\prime}, \alpha \in \Phi$, span $H$. Thus we can find a subset which forms a basis of $H$. Let $h_{\alpha_{1}}^{\prime}, \ldots, h_{\alpha_{l}}^{\prime}$ form a basis of $H$.

Proposition 5.1 Let $\alpha \in \Phi$. Then $h_{\alpha}^{\prime}=\sum_{i=1}^{l} \mu_{i} h_{\alpha_{i}}^{\prime}$ where each $\mu_{i}$ lies in $\mathbb{Q}$.

Proof. We know that $h_{\alpha}^{\prime}=\sum_{i=1}^{l} \mu_{i} h_{\alpha_{i}}^{\prime}$ for uniquely determined elements $\mu_{i} \in \mathbb{C}$. Let $\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle=\xi_{i j}$. Then $\xi_{i j} \in \mathbb{Q}$ by Proposition 4.24. We consider the system of equations:

$$
\begin{aligned}
& \left\langle h_{\alpha}^{\prime}, h_{\alpha_{1}}^{\prime}\right\rangle=\mu_{1} \xi_{11}+\mu_{2} \xi_{21}+\cdots+\mu_{l} \xi_{l 1} \\
& \left\langle h_{\alpha}^{\prime}, h_{\alpha_{2}}^{\prime}\right\rangle=\mu_{1} \xi_{12}+\mu_{2} \xi_{22}+\cdots+\mu_{l} \xi_{l 2} \\
& \quad \vdots \\
& \left\langle h_{\alpha}^{\prime}, h_{\alpha_{l}}^{\prime}\right\rangle=\mu_{1} \xi_{1 l}+\mu_{2} \xi_{2 l}+\cdots+\mu_{l} \xi_{l l} .
\end{aligned}
$$

This is a system of $l$ equations in $l$ variables $\mu_{1}, \ldots, \mu_{l}$. Now $\operatorname{det}\left(\xi_{i j}\right) \neq 0$ since the Killing form on $L$ is non-degenerate on restriction to $H$, by Proposition 4.14. Thus we may solve this system of equations for $\mu_{1}, \ldots, \mu_{l}$ by Cramer's rule. Since $\left\langle h_{\alpha}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle \in \mathbb{Q}$ and all $\xi_{i j} \in \mathbb{Q}$ we deduce that $\mu_{i} \in \mathbb{Q}$ for $i=1, \ldots, l$.

We denote by $H_{\mathbb{Q}}$ the set of all elements of form $\sum_{i=1}^{l} \mu_{i} h_{\alpha_{i}}^{\prime}$ for $\mu_{i} \in \mathbb{Q}$ and $H_{\mathbb{R}}$ the set of all such elements with $\mu_{i} \in \mathbb{R}$. Proposition 5.1 shows that $H_{\mathbb{Q}}$ and $H_{\mathbb{R}}$ are independent of the choice of basis $h_{\alpha_{i}}^{\prime}$. Also $H_{\mathbb{Q}}$ is the set of all
rational linear combinations of the $h_{\alpha}^{\prime}, \alpha \in \Phi$, and $H_{\mathbb{R}}$ is the set of all real linear combinations of such elements.

We show next that the Killing form of $L$ behaves in a favourable manner when restricted to $H_{\mathbb{R}}$.

Proposition 5.2 Let $x \in H_{\mathbb{R}}$. Then $\langle x, x\rangle \in \mathbb{R}$ and $\langle x, x\rangle \geq 0$. If $\langle x, x\rangle=0$ then $x=0$.

Proof. Let $x=\sum_{i=1}^{l} \mu_{i} h_{\alpha_{i}}^{\prime}$. Then we have

$$
\begin{aligned}
\langle x, x\rangle & =\sum_{i=1}^{l} \sum_{j=1}^{l} \mu_{i} \mu_{j}\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle \\
& =\sum_{i} \sum_{j} \mu_{i} \mu_{j} \operatorname{tr}\left(\operatorname{ad} h_{\alpha_{i}}^{\prime} \operatorname{ad} h_{\alpha_{j}}^{\prime}\right) \\
& =\sum_{i} \sum_{j} \mu_{i} \mu_{j} \sum_{\lambda \in \Phi} \lambda\left(h_{\alpha_{i}}^{\prime}\right) \lambda\left(h_{\alpha_{j}}^{\prime}\right) \\
& =\sum_{\lambda \in \Phi} \sum_{i} \sum_{j} \mu_{i} \mu_{j} \lambda\left(h_{\alpha_{i}}^{\prime}\right) \lambda\left(h_{\alpha_{j}}^{\prime}\right) \\
& =\sum_{\lambda \in \Phi}\left(\sum_{i} \mu_{i} \lambda\left(h_{\alpha_{i}}^{\prime}\right)\right)^{2} .
\end{aligned}
$$

Now $\lambda\left(h_{\alpha_{i}}^{\prime}\right)=\left\langle h_{\lambda}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle \in \mathbb{Q}$ by Proposition 4.24. Thus we have $\langle x, x\rangle \in \mathbb{R}$, and also $\langle x, x\rangle \geq 0$.

Suppose that $\langle x, x\rangle=0$. Then we have $\sum_{i} \mu_{i} \lambda\left(h_{\alpha_{i}}^{\prime}\right)=0$ for all $\lambda \in \Phi$. In particular $\sum_{i} \mu_{i} \alpha_{j}\left(h_{\alpha_{i}}^{\prime}\right)=0$ for $j=1, \ldots, l$. This gives $\sum_{i} \mu_{i}\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle=0$, that is $\sum_{i} \mu_{i} \xi_{i j}=0$. Since the matrix $\left(\xi_{i j}\right)$ is non-singular we deduce that $\mu_{i}=0$ for all $i$. Thus $x=0$.

This proposition shows that the Killing form restricted to $H_{\mathbb{R}}$ is a map $H_{\mathbb{R}} \times H_{\mathbb{R}} \rightarrow \mathbb{R}$ which is a symmetric positive definite bilinear form. The vector space $H_{\mathbb{R}}$ endowed with this positive definite form is a Euclidean space. This Euclidean space contains all vectors $h_{\alpha}^{\prime}$ for $\alpha \in \Phi$.

We recall from Lemma 4.16 that we have an isomorphism $h \rightarrow h^{*}$ from $H$ to $H^{*}$ given by $h^{*}(x)=\langle h, x\rangle$. We define $H_{\mathbb{R}}^{*}$ to be the image of $H_{\mathbb{R}}$ under this isomorphism. $H_{\mathbb{R}}^{*}$ is the real subspace of $H^{*}$ spanned by $\Phi$. We may also define a symmetric positive definite bilinear form on $H_{\mathbb{R}}^{*}$ by

$$
\left\langle h_{1}^{*}, h_{2}^{*}\right\rangle=\left\langle h_{1}, h_{2}\right\rangle \in \mathbb{R} .
$$

Thus $H_{\mathbb{R}}^{*}$ becomes a Euclidean space containing all the roots $\alpha \in \Phi$. We shall investigate the configuration formed by the roots in the Euclidean space $H_{\mathbb{R}}^{*}$. We shall, for the time being, write $V=H_{\mathbb{R}}^{*}$.

A total ordering on $V$ is a relation $<$ on $V$ satisfying the following axioms.
(i) $\lambda<\mu$ and $\mu<\nu$ implies $\lambda<\nu$.
(ii) For each pair of elements $\lambda, \mu \in V$ just one of the conditions $\lambda<\mu$, $\lambda=\mu, \mu<\lambda$ holds.
(iii) If $\lambda<\mu$ then $\lambda+\nu<\mu+\nu$.
(iv) If $\lambda<\mu$ and $\xi \in \mathbb{R}$ with $\xi>0$ then $\xi \lambda<\xi \mu$, and if $\xi<0$ then $\xi \mu<\xi \lambda$.

Every real vector space has such total orderings. If $v_{1}, \ldots, v_{l}$ are a basis of $V$ and $\lambda=\sum \lambda_{i} v_{i}, \mu=\sum \mu_{i} v_{i}$ with $\lambda \neq \mu$ then we may define $\lambda<\mu$ if the first non-zero coefficient $\mu_{i}-\lambda_{i}$ is positive. This gives us a total ordering on $V$.

A positive system $\Phi^{+} \subset \Phi$ is the set of all roots $\alpha \in \Phi$ satisfying $0<\alpha$ for some total ordering on $V$. Given such a positive system $\Phi^{+}$we define the fundamental system $\Pi \subset \Phi^{+}$as follows: $\alpha \in \Pi$ if and only if $\alpha \in \Phi^{+}$ and $\alpha$ cannot be expressed as the sum of two elements of $\Phi^{+} . \Phi^{-}$is the corresponding set of negative roots.

Proposition 5.3 Every root in $\Phi^{+}$is a sum of roots in $\Pi$.

Proof. Let $\alpha \in \Phi^{+}$. Then either $\alpha \in \Pi$ or $\alpha=\beta+\gamma$ where $\beta, \gamma \in \Phi^{+}$and $\beta<\alpha, \gamma<\alpha$. We continue this process, which must eventually terminate since $\Phi^{+}$is finite. Thus $\alpha$ is a sum of elements of $\Pi$.

Proposition 5.4 Let $\alpha, \beta \in \Pi$ with $\alpha \neq \beta$. Then $\langle\alpha, \beta\rangle \leq 0$.

Proof. We first observe that $\alpha-\beta \notin \Phi$. For if $\alpha-\beta \in \Phi$ we would have either $\alpha-\beta \in \Phi^{+}$or $\beta-\alpha \in \Phi^{+}$. If $\alpha-\beta \in \Phi^{+}$then $\alpha=(\alpha-\beta)+\beta$ which contra$\operatorname{dicts} \alpha \in \Pi$. If $\beta-\alpha \in \Phi^{+}$then $\beta=(\beta-\alpha)+\alpha$ which contradicts $\beta \in \Pi$. Hence $\alpha-\beta \notin \Phi$. We now consider the $\alpha$-chain of roots through $\beta$. This has form

$$
\beta, \alpha+\beta, \ldots, q \alpha+\beta
$$

since $-\alpha+\beta \notin \Phi$. By Proposition 4.22 we deduce

$$
2 \frac{\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}=-q
$$

However, $\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle>0$, hence $\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle \leq 0$. It follows that $\langle\alpha, \beta\rangle \leq 0$.
Thus any two distinct roots in the fundamental system $\Pi$ are inclined at an obtuse angle.

Our next result shows the importance of the concept of a fundamental system of roots.

Theorem 5.5 A fundamental system $\Pi$ forms a basis of $V=H_{\mathbb{R}}^{*}$.

Proof. We first show that $\Pi$ spans $V$. We know by Proposition 4.17 that $\Phi$ spans $V$. Since $\alpha \in \Phi$ if and only if $-\alpha \in \Phi$ we see that $\Phi^{+}$spans $V$. By Proposition 5.3 we deduce that $\Pi$ spans $V$.

We show now that the set $\Pi$ is linearly independent. Suppose this were false. Then there would exist a non-trivial linear combination of the roots $\alpha_{i} \in \Pi$ equal to zero. We take all the terms with positive coefficient to one side of this relation. Thus we have

$$
\mu_{i_{1}} \alpha_{i_{1}}+\cdots+\mu_{i_{r}} \alpha_{i_{r}}=\mu_{j_{1}} \alpha_{j_{1}}+\cdots+\mu_{j_{s}} \alpha_{j_{s}}
$$

where $\mu_{i_{1}}, \ldots, \mu_{i_{r}}, \mu_{j_{1}}, \ldots, \mu_{j_{s}}>0$ and $\alpha_{i_{1}}, \ldots, \alpha_{i_{r}}, \alpha_{j_{1}}, \ldots, \alpha_{j_{s}}$ are distinct elements of $\Pi$. We write

$$
v=\mu_{i_{1}} \alpha_{i_{1}}+\cdots+\mu_{i_{r}} \alpha_{i_{r}}=\mu_{j_{1}} \alpha_{j_{1}}+\cdots+\mu_{j_{s}} \alpha_{j_{s}} .
$$

Then we have $\langle v, v\rangle=\left\langle\mu_{i_{1}} \alpha_{i_{1}}+\cdots+\mu_{i_{r}} \alpha_{i_{r}}, \mu_{j_{1}} \alpha_{j_{1}}+\cdots+\mu_{j_{s}} \alpha_{j_{s}}\right\rangle$. We deduce $\langle v, v\rangle \leq 0$ by Proposition 5.4. Since the form is positive definite this implies that $v=0$. However, $0<v$ since we have $0<\alpha_{i}$ for all $\alpha_{i} \in \Pi$ and $\mu_{i}>0$. This gives a contradiction. Thus $\Pi$ is linearly independent.

We see in particular that $|\Pi|=l=\operatorname{dim} H$. Thus the number of roots in a fundamental system is equal to the rank of the Lie algebra $L$.

Corollary 5.6 Let $\Pi$ be a fundamental system of roots. Then each $\alpha \in \Phi$ can be expressed in the form $\alpha=\sum n_{i} \alpha_{i}$ where $\alpha_{i} \in \Pi, n_{i} \in \mathbb{Z}$ and either $n_{i} \geq 0$ for all $i$ or $n_{i} \leq 0$ for all $i$.

Proof. The roots $\alpha \in \Phi^{+}$have all $n_{i} \geq 0$ and the roots $\alpha \in \Phi^{-}$have all $n_{i} \leq 0$.

### 5.2 The Weyl group

Inside the root system $\Phi$ a positive system $\Phi^{+}$can be chosen in many different ways. However, we shall show that any two positive systems in $\Phi$ can be transformed into one another by an element of a certain finite group $W$ which acts on $\Phi$.

For each $\alpha \in \Phi$ we define a linear map $s_{\alpha}: V \rightarrow V$ by

$$
s_{\alpha}(x)=x-2 \frac{\langle\alpha, x\rangle}{\langle\alpha, \alpha\rangle} \alpha \quad \text { for all } x \in V
$$

As before, $V=H_{\mathbb{R}}^{*}$. This map $s_{\alpha}$ satisfies

$$
\begin{aligned}
& s_{\alpha}(\alpha)=-\alpha \\
& s_{\alpha}(x)=x \quad \text { if }\langle\alpha, x\rangle=0 .
\end{aligned}
$$

There is a unique linear map satisfying these conditions - the reflection in the hyperplane of $V$ orthogonal to $\alpha$. Thus $s_{\alpha}$ is this reflection.

The group $W$ of all non-singular linear maps on $V$ generated by the $s_{\alpha}$ for all $\alpha \in \Phi$ is called the Weyl group. This group plays an important role in the Lie theory. It is a group of isometries of $V$, that is we have

$$
\langle w x, w y\rangle=\langle x, y\rangle \quad \text { for all } x, y \in V
$$

Proposition 5.7 $W$ permutes the roots. Thus if $\alpha \in \Phi$ and $w \in W$ then $w(\alpha) \in$ $\Phi$.

Proof. It is sufficient to show that $s_{\alpha}(\beta) \in \Phi$ for all $\alpha, \beta \in \Phi$ since the elements $s_{\alpha}$ generate $W$. If $\beta=\alpha$ or $-\alpha$ this is clear. Thus suppose $\beta \neq \pm \alpha$. Let the $\alpha$-chain of roots through $\beta$ be

$$
-p \alpha+\beta, \ldots, \beta, \ldots, q \alpha+\beta
$$

Then we have

$$
s_{\alpha}(\beta)=\beta-2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle} \alpha=\beta-(p-q) \alpha
$$

by Proposition 4.22. Now $\beta-(p-q) \alpha$ is one of the roots in the $\alpha$-chain through $\beta$. Thus $s_{\alpha}(\beta) \in \Phi$.

In fact we observe that $s_{\alpha}$ inverts the above $\alpha$-chain of roots. In particular we have

$$
s_{\alpha}(q \alpha+\beta)=-p \alpha+\beta, \quad s_{\alpha}(-p \alpha+\beta)=q \alpha+\beta
$$

Proposition 5.8 The Weyl group $W$ is finite.

Proof. $W$ permutes $\Phi$ and $\Phi$ is finite. If two elements of $W$ induce the same permutation of $\Phi$ they must be equal, since $\Phi$ spans $V$. Since there are only finitely many permutations of $\Phi, W$ must be finite.

Now suppose that $\Phi^{+}$is a positive system in $\Phi$ and that $\Pi$ is the corresponding fundamental system.

Lemma 5.9 Let $\alpha \in \Pi$. If $\beta \in \Phi^{+}$and $\beta \neq \alpha$ then $s_{\alpha}(\beta) \in \Phi^{+}$.

Proof. We can express $\beta$ in the form

$$
\beta=\sum_{i} n_{i} \alpha_{i} \quad \alpha_{i} \in \Pi, \quad n_{i} \in \mathbb{Z}, \quad n_{i} \geq 0
$$

by Corollary 5.6. Since $\beta \neq \alpha$ there must be some $n_{i} \neq 0$ with $\alpha_{i} \neq \alpha$. We then consider

$$
s_{\alpha}(\beta)=\beta-2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle} \alpha
$$

and express this as a linear combination of the elements of $\Pi$. The coefficient of $\alpha_{i}$ in $s_{\alpha}(\beta)$ remains $n_{i}$. Since $n_{i}>0$ we deduce from Corollary 5.6 that $s_{\alpha}(\beta) \in \Phi^{+}$.

Theorem 5.10 Let $\Phi_{1}^{+}$, $\Phi_{2}^{+}$be two positive systems in $\Phi$. Then there exists $w \in W$ such that $w\left(\Phi_{1}^{+}\right)=\Phi_{2}^{+}$.

Proof. Let $m=\left|\Phi_{1}^{+} \cap \Phi_{2}^{-}\right|$. We shall use induction on $m$. If $m=0$ we have $\Phi_{1}^{+}=\Phi_{2}^{+}$and so $w=1$ has the required property. Thus we may assume $m>0$.

Let $\Pi_{1}$ be the fundamental system in $\Phi_{1}^{+}$. We cannot have $\Pi_{1} \subset \Phi_{2}^{+}$as this would imply $\Phi_{1}^{+} \subset \Phi_{2}^{+}$, contrary to $m>0$. Thus there exists $\alpha \in \Pi_{1} \cap \Phi_{2}^{-}$.

We consider $s_{\alpha}\left(\Phi_{1}^{+}\right)$. This is also a positive system in $\Phi$. By Lemma 5.9 $s_{\alpha}\left(\Phi_{1}^{+}\right)$contains all roots in $\Phi_{1}^{+}$except $\alpha$, together with $-\alpha$. Thus we have

$$
\left|s_{\alpha}\left(\Phi_{1}^{+}\right) \cap \Phi_{2}^{-}\right|=m-1
$$

By induction there exists $w^{\prime} \in W$ such that $w^{\prime} s_{\alpha}\left(\Phi_{1}^{+}\right)=\Phi_{2}^{+}$. Let $w=w^{\prime} s_{\alpha}$. Then $w\left(\Phi_{1}^{+}\right)=\Phi_{2}^{+}$as required.

Corollary 5.11 Let $\Pi_{1}, \Pi_{2}$ be two fundamental systems in $\Phi$. Then there exists $w \in W$ such that $w\left(\Pi_{1}\right)=\Pi_{2}$.

Proof. Let $\Phi_{1}^{+}, \Phi_{2}^{+}$be positive systems containing $\Pi_{1}, \Pi_{2}$ respectively. Let $\Phi_{2}^{+}=w\left(\Phi_{1}^{+}\right)$. Then $w\left(\Pi_{1}\right)$ is a fundamental system contained in $\Phi_{2}^{+}$, so $w\left(\Pi_{1}\right)=\Pi_{2}$.

Proposition 5.12 Let $\Pi$ be a fundamental system in $\Phi$. Then for each $\alpha \in \Phi$ there exist $\alpha_{i} \in \Pi$ and $w \in W$ with $\alpha=w\left(\alpha_{i}\right)$.

Proof. Let $\Phi^{+}$be the positive system with fundamental system $\Pi$. First suppose $\alpha \in \Phi^{+}$. Then we have

$$
\alpha=\sum_{i} n_{i} \alpha_{i} \quad \alpha_{i} \in \Pi, \quad n_{i} \in \mathbb{Z}, \quad n_{i} \geq 0
$$

by Corollary 5.6. We define the height of $\alpha$ by

$$
\text { ht } \alpha=\sum_{i} n_{i}
$$

We shall argue by induction on ht $\alpha$. If ht $\alpha=1$ then $\alpha=\alpha_{i}$ for some $i$ and $\alpha \in \Pi$. The result is obvious in this case. Thus suppose ht $\alpha>1$. Then we have $n_{i}>0$ for at least two values of $i$ by Proposition 4.21. Now

$$
\langle\alpha, \alpha\rangle=\sum_{i} n_{i}\left\langle\alpha, \alpha_{i}\right\rangle
$$

Since $\langle\alpha, \alpha\rangle>0$ and each $n_{i} \geq 0$ there exist $\alpha_{i} \in \Pi$ with $\left\langle\alpha, \alpha_{i}\right\rangle>0$. Let $s_{i}(\alpha)=\beta$. Then $\beta \in \Phi$ and

$$
\beta=\alpha-2 \frac{\left\langle\alpha_{i}, \alpha\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} \alpha_{i}
$$

Since $\left\langle\alpha_{i}, \alpha\right\rangle>0$ we see that ht $\beta<$ ht $\alpha$. On the other hand $\beta \in \Phi^{+}$since only one coefficient $n_{i}$ is changed in passing from $\alpha$ to $\beta$, thus at least one coefficient remains positive in $\beta$. By Corollary 5.6 this is sufficient to show that $\beta \in \Phi^{+}$. By induction there exist $\alpha_{j} \in \Pi$ and $w^{\prime} \in W$ such that $\beta=w^{\prime}\left(\alpha_{j}\right)$. Then

$$
\alpha=s_{i}(\beta)=s_{i} w^{\prime}\left(\alpha_{j}\right)
$$

as required.
Finally we suppose that $\alpha \in \Phi^{-}$. Then $\alpha=s_{\alpha}(-\alpha)$ and $-\alpha \in \Phi^{+}$. Thus $-\alpha=w^{\prime}\left(\alpha_{i}\right)$ for some $w^{\prime} \in W, \alpha_{i} \in \Pi$. Hence $\alpha=s_{\alpha} w^{\prime}\left(\alpha_{i}\right)$ as required.

Thus each root is the image of some fundamental root under an element of the Weyl group.

We show next that $W$ is generated by the reflections corresponding to roots in a given fundamental system.

Theorem 5.13 Let $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$ be a fundamental system in $\Phi$. Then the corresponding fundamental reflections $s_{\alpha_{1}}, \ldots, s_{\alpha_{l}}$ generate $W$.

Proof. Let $W_{0}$ be the subgroup of $W$ generated by $s_{\alpha_{1}}, \ldots, s_{\alpha_{l}}$. Since the $s_{\alpha}$ generate $W$ for all $\alpha \in \Phi$ it is sufficient to show that each $s_{\alpha}$ lies in $W_{0}$. We may assume $\alpha \in \Phi^{+}$since $s_{\alpha}=s_{-\alpha}$. Now the proof of Proposition 5.12
shows that $\alpha=w\left(\alpha_{i}\right)$ for some $\alpha_{i} \in \Pi$ and some $w \in W_{0}$. We consider the element $w s_{\alpha_{i}} w^{-1} \in W_{0}$. We have

$$
w s_{\alpha_{i}} w^{-1}(\alpha)=w s_{\alpha_{i}}\left(\alpha_{i}\right)=w\left(-\alpha_{i}\right)=-\alpha
$$

We shall also show $w s_{\alpha_{i}} w^{-1}(x)=x$ if $\langle\alpha, x\rangle=0$. For $\langle\alpha, x\rangle=0$ implies $\left\langle w^{-1}(\alpha), w^{-1}(x)\right\rangle=0$, that is $\left\langle\alpha_{i}, w^{-1}(x)\right\rangle=0$. This gives $s_{\alpha_{i}} w^{-1}(x)=$ $w^{-1}(x)$, i.e. $w s_{\alpha_{i}} w^{-1}(x)=x$. Thus $w s_{\alpha_{i}} w^{-1}$ is the reflection in the hyperplane orthogonal to $\alpha$, that is $w s_{\alpha_{i}} w^{-1}=s_{\alpha}$. This shows that $s_{\alpha} \in W_{0}$. Hence $W_{0}=W$.

We now wish to obtain further information about the way in which the Weyl group $W$ is generated by a set of its fundamental reflections. As before we let $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$ be a fundamental system of roots and consider the corresponding set of fundamental reflections. For simplicity we write

$$
s_{1}=s_{\alpha_{1}}, \quad s_{2}=s_{\alpha_{2}}, \quad \ldots, \quad s_{l}=s_{\alpha_{l}} .
$$

Then each element of $W$ can be expressed as a product of elements $s_{i}$. (We do not need to introduce inverses since $s_{i}^{-1}=s_{i}$.) For each $w \in W$ we define $l(w)$ to be the minimal value of $m$ such that $w$ can be expressed as a product of $m$ fundamental reflections $s_{i} . l(w)$ is called the length of $w$. It is clear that $l(1)=0$ and $l\left(s_{i}\right)=1$. An expression of $w$ as a product of fundamental reflections $s_{i}$ with $l(w)$ terms is called a reduced expression for $w$.

We shall relate $l(w)$ to another integer $n(w)$. We recall that each element $w \in W$ permutes the elements of $\Phi$. We define $n(w)$ to be the number of roots $\alpha \in \Phi^{+}$for which $w(\alpha) \in \Phi^{-}$. Thus $n(w)$ is the number of positive roots made negative by $w$. We aim to show that $l(w)=n(w)$.

Proposition $5.14 n(w) \leq l(w)$ for all $w \in W$.

Proof. We shall first compare $n(w)$ with $n\left(w s_{i}\right)$. We recall from Lemma 5.9 that $s_{i}$ transforms $\alpha_{i}$ to $-\alpha_{i}$ and all positive roots other than $\alpha_{i}$ to positive roots. It follows that

$$
n\left(w s_{i}\right)=n(w) \pm 1
$$

In order to determine the sign we consider the effect of $w$ and $w s_{i}$ on $\alpha_{i}$. If $w\left(\alpha_{i}\right) \in \Phi^{+}$then $w$ transforms $\alpha_{i}$ to a positive root and $w s_{i}$ transforms $\alpha_{i}$ to a negative root. Hence $n\left(w s_{i}\right)=n(w)+1$. On the other hand if $w\left(\alpha_{i}\right) \in \Phi^{-}$ then we get the reverse situation and $n\left(w s_{i}\right)=n(w)-1$.

Now let us take a reduced expression

$$
w=s_{i_{1}} s_{i_{2}} \ldots s_{i_{r}} \quad r=l(w)
$$

Then we have

$$
n(w) \leq n\left(s_{i_{1}} \ldots s_{i_{r-1}}\right)+1 \leq n\left(s_{i_{1}} \ldots s_{i_{r-2}}\right)+2 \leq \cdots \leq r
$$

Thus $n(w) \leq l(w)$ as required.
In order to prove the converse result $l(w) \leq n(w)$ we shall first prove a result called the deletion condition which is important in its own right.

Theorem 5.15 Let $w=s_{i_{1}} \ldots s_{i_{r}}$ be any expression of $w \in W$ as a product of fundamental reflections. Suppose $n(w)<r$. Then there exist integers $j, k$ with $1 \leq j<k \leq r$ such that

$$
w=s_{i_{1}} \ldots \hat{s}_{i_{j}} \ldots \hat{s}_{i_{k}} \ldots s_{i_{r}}
$$

where ^ denotes omission.

Proof. We recall from the proof of Proposition 5.14 that, for all $w \in W$, $n\left(w s_{i}\right)=n(w) \pm 1$. Consider the given expression

$$
w=s_{i_{1}} \ldots s_{i_{r}}
$$

Since $n(w)<r$ there exists $k$ with $1<k \leq r$ such that

$$
n\left(s_{i_{1}} \ldots s_{i_{k}}\right)=n\left(s_{i_{1}} \ldots s_{i_{k-1}}\right)-1
$$

This implies $s_{i_{1}} \ldots s_{i_{k-1}}\left(\alpha_{i_{k}}\right) \in \Phi^{-}$as in the proof of Proposition 5.14. Since $\alpha_{i_{k}} \in \Phi^{+}$there exists $j$ with $1 \leq j<k$ such that

$$
\begin{aligned}
& s_{i_{j+1}} \ldots s_{i_{k-1}}\left(\alpha_{i_{k}}\right) \in \Phi^{+} \\
& s_{i_{j}} s_{i_{j+1}} \ldots s_{i_{k-1}}\left(\alpha_{i_{k}}\right) \in \Phi^{-} .
\end{aligned}
$$

By Lemma $5.9 s_{i_{j}}$ transforms only one positive root into a negative root, namely $\alpha_{i_{j}}$. Thus we have

$$
s_{i_{j+1}} \ldots s_{i_{k-1}}\left(\alpha_{i_{k}}\right)=\alpha_{i_{j}}
$$

It follows that the reflections $s_{i_{k}}, s_{i_{j}}$ associated with the roots $\alpha_{i_{k}}, \alpha_{i_{j}}$ are related by

$$
s_{i_{j}}=s_{i_{j+1}} \ldots s_{i_{k-1}} s_{i_{k}} s_{i_{k-1}} \ldots s_{i_{j+1}}
$$

This implies

$$
s_{i_{j}} s_{i_{j+1}} \ldots s_{i_{k-1}}=s_{i_{j+1}} \ldots s_{i_{k-1}} s_{i_{k}}
$$

Thus we have

$$
s_{i_{1}} \ldots s_{i_{r}}=s_{i_{1}} \ldots s_{i_{j-1}} s_{i_{j+1}} \ldots s_{i_{k-1}} s_{i_{k+1}} \ldots s_{i_{r}}
$$

and so $w=s_{i_{1}} \ldots \hat{s}_{i_{j}} \ldots \hat{s}_{i_{k}} \ldots s_{i_{r}}$ as required.
Corollary $5.16 n(w)=l(w)$.
Proof. We know from Proposition 5.14 that $n(w) \leq l(w)$. Suppose if possible that $n(w)<l(w)$. Let $w=s_{i_{1}} \ldots s_{i_{r}}$ be a reduced expression, thus $r=l(w)$. Since $n(w)<r$ we may apply Theorem 5.15 to show that $w$ is a product of $r-2$ fundamental reflections. This contradicts the definition of $l(w)$.

Thus the length of $w$ is equal to the number of positive roots made negative by $w$.

Proposition 5.17 (a) The maximal length of any element of $W$ is $\left|\Phi^{+}\right|$.
(b) $W$ has a unique element $w_{0}$ with $l\left(w_{0}\right)=\left|\Phi^{+}\right|$.
(c) $w_{0}\left(\Phi^{+}\right)=\Phi^{-}$
(d) $w_{0}^{2}=1$.

Proof. Since $l(w)=n(w)$ we have $l(w) \leq\left|\Phi^{+}\right|$. For each fundamental system $\Pi \subset \Phi,-\Pi$ is also a fundamental system, coming from the opposite total ordering. Thus by Corollary 5.11 there exists $w_{0} \in W$ with $w_{0}(\Pi)=-\Pi$. Hence $w_{0}\left(\Phi^{+}\right)=\Phi^{-}$and $n\left(w_{0}\right)=\left|\Phi^{+}\right|$. Thus $l\left(w_{0}\right)=\left|\Phi^{+}\right|$also and $w_{0}$ is an element of $W$ of maximal length.

Now let $w_{0}^{\prime} \in W$ also have $l\left(w_{0}^{\prime}\right)=\left|\Phi^{+}\right|$. Then $n\left(w_{0}^{\prime}\right)=\left|\Phi^{+}\right|$and so $w_{0}^{\prime}\left(\Phi^{+}\right)=\Phi^{-}$. Let $w=\left(w_{0}^{\prime}\right)^{-1} w_{0}$. Then $w\left(\Phi^{+}\right)=\Phi^{+}$and so $n(w)=0$. Hence $l(w)=0$ and so $w=1$. Thus $w_{0}^{\prime}=w_{0}$ and the element $w_{0}$ of maximal length is unique.

Finally we have $w_{0}^{2}\left(\Phi^{+}\right)=\Phi^{+}$and so $n\left(w_{0}^{2}\right)=0$. Hence $l\left(w_{0}^{2}\right)=0$ and $w_{0}^{2}=1$.

### 5.3 Generators and relations for the Weyl group

In this section we shall give a description of the Weyl group $W$ by means of generators and relations. Let the order of the element $s_{i} s_{j} \in W$ be $m_{i j}$ when $i \neq j$.

Theorem 5.18 W is isomorphic to the abstract group given by generators and relations:

$$
\left\langle s_{1}, \ldots, s_{l} ; s_{i}^{2}=1, \quad\left(s_{i} s_{j}\right)^{m_{i j}}=1 \quad \text { for } i \neq j\right\rangle
$$

A group defined by generators and relations of this form is called a Coxeter group. Thus the theorem asserts that the Weyl group is a Coxeter group.

Proof. Since $W$ is generated by $s_{1}, \ldots, s_{l}$ and the relations $s_{i}^{2}=1$ and $\left(s_{i} s_{j}\right)^{m_{i j}}=1$ hold in $W$ it is sufficient to show that every relation

$$
s_{i_{1}} \ldots s_{i_{r}}=1
$$

in $W$ is a consequence of the defining relations. Now each $s_{i}$ is a reflection, thus det $s_{i}=-1$. Hence $\operatorname{det}\left(s_{i_{1}} \ldots s_{i_{r}}\right)=(-1)^{r}$. If $s_{i_{1}} \ldots s_{i_{r}}=1$ we deduce that $r$ must be even. Let $r=2 q$. We shall show that

$$
s_{i_{1}} \ldots s_{i_{2 q}}=1
$$

is a consequence of the defining relations, by induction on $q$. If $q=1$ the relation is $s_{i_{1}} s_{i_{2}}=1$, hence $s_{i_{2}}=s_{i_{1}}^{-1}=s_{i_{1}}$. Our relation is thus $s_{i_{1}}^{2}=1$, which is one of the defining relations.

We may therefore assume inductively that all relations in $W$ of length less than $2 q$ are consequences of the defining relations.

Now the given relation can be written

$$
s_{i_{1}} \ldots s_{i_{q}} s_{i_{q+1}}=s_{i_{2 q}} \ldots s_{i_{q+2}}
$$

Thus $l\left(s_{i_{1}} \ldots s_{i_{q}} s_{i_{q+1}}\right)<q+1$. Hence, by the deletion condition Theorem 5.15, we have

$$
s_{i_{1}} \ldots s_{i_{q+1}}=s_{i_{1}} \ldots \hat{s}_{i_{j}} \ldots \hat{s}_{i_{k}} \ldots s_{i_{q+1}}
$$

for certain $j, k$ with $1 \leq j<k \leq q+1$. Now unless $j=1$ and $k=q+1$ this is a consequence of a relation with fewer than $2 q$ terms. It can therefore be deduced from the defining relations. The relation

$$
s_{i_{1}} \ldots \hat{s}_{i_{j}} \ldots \hat{s}_{i_{k}} \ldots s_{i_{q+1}}=s_{i_{2 q}} \ldots s_{i_{q+2}}
$$

has $2 q-2$ terms, so is also a consequence of the defining relations. Thus the given relation

$$
s_{i_{1}} \ldots s_{i_{q+1}}=s_{i_{2 q}} \ldots s_{i_{q+2}}
$$

will be a consequence of the defining relations, unless we have $j=1$ and $k=q+1$.

We may therefore assume that $j=1$ and $k=q+1$. Thus we have

$$
s_{i_{1}} \ldots s_{i_{q+1}}=s_{i_{2}} \ldots s_{i_{q}}
$$

that is

$$
s_{i_{1}} \ldots s_{i_{q}}=s_{i_{2}} \ldots s_{i_{q+1}} .
$$

We now write the original relation

$$
s_{i_{1}} \ldots s_{i_{2 q}}=1
$$

in the alternative form

$$
s_{i_{2}} \ldots s_{i_{2 q}} s_{i_{1}}=1
$$

In exactly the same way this relation will be a consequence of the defining relations unless

$$
s_{i_{2}} \ldots s_{i_{q+1}}=s_{i_{3}} \ldots s_{i_{q+2}} .
$$

If this relation is a consequence of the defining relations then the relation

$$
s_{i_{2}} \ldots s_{i_{2 q}} s_{i_{1}}=1
$$

will also be a consequence of the defining relations, by the above argument, and we are done.

Now $s_{i_{2}} \ldots s_{i_{q+1}}=s_{i_{3}} \ldots s_{i_{q+2}}$ is equivalent to

$$
s_{i_{3}} s_{i_{2}} s_{i_{3}} \ldots s_{i_{q}} s_{i_{q+1}} s_{i_{q+2}} s_{i_{q+1}} \ldots s_{i_{4}}=1
$$

and this will be a consequence of the defining relations unless

$$
s_{i_{3}} s_{i_{2}} s_{i_{3}} \ldots s_{i_{q}}=s_{i_{2}} s_{i_{3}} \ldots s_{i_{q}} s_{i_{q+1}}
$$

We may therefore assume this to be true. But we also have

$$
s_{i_{1}} s_{i_{2}} s_{i_{3}} \ldots s_{i_{q}}=s_{i_{2}} s_{i_{3}} \ldots s_{i_{q}} s_{i_{q+1}}
$$

and so $s_{i_{1}}=s_{i_{3}}$. Hence the given relation

$$
s_{i_{1}} \ldots s_{i_{2 q}}=1
$$

will be a consequence of the defining relations unless $s_{i_{1}}=s_{i_{3}}$.
However, the given relation can be written in the equivalent forms

$$
\begin{aligned}
& s_{i_{2}} \ldots s_{i_{2 q}} s_{i_{1}}=1 \\
& s_{i_{3}} \ldots s_{i_{2 q}} s_{i_{1}} s_{i_{2}}=1
\end{aligned}
$$

and so on. Thus this relation will be a consequence of the defining relations unless we have

$$
\begin{aligned}
& s_{i_{1}}=s_{i_{3}}=s_{i_{5}}=\ldots=s_{i_{2 q-1}} \\
& s_{i_{2}}=s_{i_{4}}=s_{i_{6}}=\ldots=s_{i_{2 q}}
\end{aligned}
$$

Thus we may assume that the given relation has form

$$
s_{i_{1}} s_{i_{2}} s_{i_{1}} s_{i_{2}} \ldots s_{i_{1}} s_{i_{2}}=1
$$

that is $\left(s_{i_{1}} s_{i_{2}}\right)^{q}=1$. Now the order of $s_{i_{1}} s_{i_{2}}$ is $m_{i_{1} i_{2}}$, hence $m_{i_{1} i_{2}}$ divides $q$. Thus the relation $\left(s_{i_{1}} s_{i_{2}}\right)^{q}=1$ is a consequence of the defining relation $\left(s_{i_{1}} s_{i_{2}}\right)^{m_{i_{1} i_{2}}}=1$. This completes the proof.

This remarkable proof, due to R. Steinberg, shows that the Weyl group W is a finite Coxeter group.

## 6

## The Cartan matrix and the Dynkin diagram

### 6.1 The Cartan matrix

We shall now investigate in more detail the geometry of the system of roots $\Phi$ in the vector space $V=H_{\mathbb{R}}^{*}$. We recall from Proposition 5.2 that $V$ is a Euclidean space with respect to the scalar product $\langle$,$\rangle . The roots \Phi$ span $V$ but are not linearly independent. Any fundamental system $\Pi \subset \Phi$ forms a basis of $V$.
We first consider the possible angles between pairs of roots $\alpha, \beta \in \Phi$ and the relative lengths of the roots $\alpha, \beta$. The angles will be taken to satisfy $0 \leq \theta \leq \pi$.

Proposition 6.1 Let $\alpha, \beta \in \Phi$ be such that $\beta \neq \pm \alpha$. Then:
(i) the angle between $\alpha, \beta$ is one of $\pi / 6, \pi / 4, \pi / 3, \pi / 2,2 \pi / 3,3 \pi / 4,5 \pi / 6$
(ii) if $\alpha, \beta$ are inclined at $\pi / 3$ or $2 \pi / 3$ then $\alpha, \beta$ have the same length
(iii) if $\alpha, \beta$ are inclined at $\pi / 4$ or $3 \pi / 4$ then the ratio of their lengths is $\sqrt{ } 2$
(iv) if $\alpha, \beta$ are inclined at $\pi / 6$ or $5 \pi / 6$ then the ratio of their lengths is $\sqrt{ } 3$.

Proof. Let $\theta$ be the angle between $\alpha, \beta$. Then we have

$$
\langle\alpha, \beta\rangle=|\alpha||\beta| \cos \theta
$$

where $|\alpha|=\sqrt{ }\langle\alpha, \alpha\rangle$. Hence

$$
\cos ^{2} \theta=\frac{\langle\alpha, \beta\rangle^{2}}{\langle\alpha, \alpha\rangle\langle\beta, \beta\rangle}=\frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle} \cdot \frac{\langle\beta, \alpha\rangle}{\langle\beta, \beta\rangle}
$$

and so

$$
4 \cos ^{2} \theta=2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle} \cdot 2 \frac{\langle\beta, \alpha\rangle}{\langle\beta, \beta\rangle} .
$$

Now we recall from Proposition 4.22 that $2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle}$ and $2 \frac{\langle\beta, \alpha\rangle}{\langle\beta, \beta\rangle}$ are integers. Hence $4 \cos ^{2} \theta \in \mathbb{Z}$. Since $0 \leq 4 \cos ^{2} \theta \leq 4$ and $\beta \neq \pm \alpha$ we have $4 \cos ^{2} \theta \in\{0,1,2,3\}$. We consider in each case the possible factorisations of $4 \cos ^{2} \theta$ into the product of two integers.

First suppose $4 \cos ^{2} \theta=0$. Then $\theta=\pi / 2$.
Next suppose $4 \cos ^{2} \theta=1$. Then $\cos \theta=\frac{1}{2}$ or $-\frac{1}{2}$, hence $\theta=\pi / 3$ or $2 \pi / 3$. The possible factorisations of $4 \cos ^{2} \theta$ are

$$
1=1 \cdot 1 \text { or } 1=-1 \cdot-1
$$

In either case we have

$$
2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle}=2 \frac{\langle\beta, \alpha\rangle}{\langle\beta, \beta\rangle}
$$

and so $\langle\alpha, \alpha\rangle=\langle\beta, \beta\rangle$ and $\alpha, \beta$ have the same length.
Next suppose $4 \cos ^{2} \theta=2$. Then $\cos \theta=1 / \sqrt{ } 2$ or $-1 / \sqrt{ } 2$, thus $\theta=\pi / 4$ or $3 \pi / 4$. The possible factorisations of $4 \cos ^{2} \theta$ are

$$
2=1 \cdot 2 \text { or } 2=-1 \cdot-2
$$

In either case, by choosing $\alpha, \beta$ in a suitable order, we have

$$
2 \frac{\langle\beta, \alpha\rangle}{\langle\beta, \beta\rangle}=2 \cdot 2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle}
$$

that is $\langle\alpha, \alpha\rangle=2\langle\beta, \beta\rangle$ and $|\alpha|=\sqrt{ } 2|\beta|$. Thus the ratio of the lengths of $\alpha, \beta$ is $\sqrt{ } 2$.

Finally suppose that $4 \cos ^{2} \theta=3$. Then $\cos \theta=\sqrt{ } 3 / 2$ or $-\sqrt{ } 3 / 2$, so $\theta=\pi / 6$ or $5 \pi / 6$. The possible factorisations of $4 \cos ^{2} \theta$ are

$$
3=1 \cdot 3 \text { or } 3=-1 \cdot-3
$$

In either case, by choosing $\alpha, \beta$ in a suitable order, we have

$$
2 \frac{\langle\beta, \alpha\rangle}{\langle\beta, \beta\rangle}=3 \cdot 2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle}
$$

that is $\langle\alpha, \alpha\rangle=3\langle\beta, \beta\rangle$ and $|\alpha|=\sqrt{ } 3|\beta|$. Thus the ratio of the lengths of $\alpha, \beta$ is $\sqrt{ } 3$.

This completes the proof. We do not obtain any information about the relative lengths of $\alpha, \beta$ in the case when $\theta=\pi / 2$.

Corollary 6.2 Let $\Pi$ be a fundamental system of roots and let $\alpha, \beta \in \Pi$ with $\beta \neq \alpha$. Then the angle between $\alpha, \beta$ is one of $\pi / 2,2 \pi / 3,3 \pi / 4,5 \pi / 6$.

Proof. This follows from Proposition 6.1 together with the fact, proved in Proposition 5.4, that the angle $\theta$ between two distinct fundamental roots satisfies $\pi / 2 \leq \theta<\pi$.

Let $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$ be a fundamental system. We incorporate the information about the angles between the $\alpha_{i}$ and their relative lengths in the form of a matrix. We define $A_{i j}$ by

$$
A_{i j}=2 \frac{\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} \quad i, j=1, \ldots, l
$$

Thus $A_{i j} \in \mathbb{Z}$. The $l \times l$ matrix $A=\left(A_{i j}\right)$ is called the Cartan matrix.
Proposition 6.3 The Cartan matrix A has the following properties.
(i) $A_{i i}=2$ for all $i$.
(ii) $A_{i j} \in\{0,-1,-2,-3\}$ if $i \neq j$.
(iii) If $A_{i j}=-2$ or -3 then $A_{j i}=-1$.
(iv) $A_{i j}=0$ if and only if $A_{j i}=0$.

Proof. Properties (i), (iv) are obvious and (ii), (iii) follow from the proof of Proposition 6.1.

If we number the fundamental roots in $\Pi$ in a different way we may well get a different Cartan matrix $A$. However, apart from this ambiguity of numbering, the Cartan matrix $A$ is uniquely determined by the semisimple Lie algebra $L$.

Proposition 6.4 The Cartan matrix of Lepends only on the numbering of the fundamental roots. It is independent of the choice of Cartan subalgebra $H$ and fundamental system $\Pi$.

Proof. The independence of the choice of Cartan subalgebra follows from the conjugacy of Cartan subalgebras, proved in Theorem 3.13.

Let $\Pi^{\prime}$ be a second fundamental system. By Corollary 5.11 there exists $w \in W$ with $w(\Pi)=\Pi^{\prime}$. Let $w\left(\alpha_{i}\right)=\alpha_{i}^{\prime}$. Since $w$ is an isometry of $V$ we have

$$
2 \frac{\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}=2 \frac{\left\langle\alpha_{i}^{\prime}, \alpha_{j}^{\prime}\right\rangle}{\left\langle\alpha_{i}^{\prime}, \alpha_{i}^{\prime}\right\rangle}
$$

Thus the Cartan matrices defined by $\Pi$ and $\Pi^{\prime}$ with respect to these labellings are the same.

The only possible $1 \times 1$ Cartan matrix is (2). We also see that any $2 \times 2$ Cartan matrix must be one of the following:

$$
\left(\begin{array}{ll}
2 & 0 \\
0 & 2
\end{array}\right)\left(\begin{array}{cc}
2 & -1 \\
-1 & 2
\end{array}\right) \quad\left(\begin{array}{cc}
2 & -1 \\
-2 & 2
\end{array}\right) \quad\left(\begin{array}{cc}
2 & -2 \\
-1 & 2
\end{array}\right) \quad\left(\begin{array}{cc}
2 & -1 \\
-3 & 2
\end{array}\right)\left(\begin{array}{cc}
2 & -3 \\
-1 & 2
\end{array}\right)
$$

The pair

$$
\left(\begin{array}{cc}
2 & -1 \\
-2 & 2
\end{array}\right) \quad\left(\begin{array}{cc}
2 & -2 \\
-1 & 2
\end{array}\right)
$$

are obtained from one another by reversing the labelling 1,2 , and so are the pair

$$
\left(\begin{array}{cc}
2 & -1 \\
-3 & 2
\end{array}\right) \quad\left(\begin{array}{cc}
2 & -3 \\
-1 & 2
\end{array}\right)
$$

### 6.2 The Dynkin diagram

In order to determine the possible $l \times l$ Cartan matrices for larger values of $l$ it is useful to introduce a graph called the Dynkin diagram. The Dynkin diagram is determined by the Cartan matrix. It is a graph with vertices labelled $1, \ldots, l$. If $i \neq j$ the vertices $i, j$ are joined by $n_{i j}$ edges, where

$$
n_{i j}=A_{i j} A_{j i}
$$

We see from Proposition 6.4 that the Dynkin diagram is uniquely determined by the semisimple Lie algebra $L$.

The Dynkin diagrams of the Cartan matrices of degrees 1 and 2 are as follows.

## Cartan matrix Dynkin diagram

$$
\begin{gathered}
(2) \\
\left(\begin{array}{ll}
2 & 0 \\
0 & 2
\end{array}\right) \\
\left(\begin{array}{cc}
2 & -1 \\
-1 & 2
\end{array}\right) \\
\left(\begin{array}{cc}
2 & -1 \\
-2 & 2
\end{array}\right)\left(\begin{array}{cc}
2 & -2 \\
-1 & 2
\end{array}\right) \\
\left(\begin{array}{cc}
2 & -1 \\
-3 & 2
\end{array}\right)\left(\begin{array}{cc}
2 & -3 \\
-1 & 2
\end{array}\right)
\end{gathered}
$$

Proposition $6.5 n_{i j} \in\{0,1,2,3\}$ for all $i \neq j$.

Proof. This follows from Proposition 6.3 and the fact that $n_{i j}=A_{i j} A_{j i}$

Thus the number of edges joining any two distinct vertices of the Dynkin diagram is either $0,1,2$ or 3 .

Now the Dynkin diagram need not be a connected graph. However, if it is disconnected it will split into connected components. If we number the vertices so that those in each connected component are numbered consecutively, the Cartan matrix will split into blocks of the form

$$
A=\left(\begin{array}{c|c|c|c}
* & O & O & O \\
\hline O & * & O & O \\
\hline O & O & * & O \\
\hline O & O & O & *
\end{array}\right) .
$$

with one diagonal block for each connected component. This diagonal block will be the Cartan matrix for the given connected component. The set $\Pi=$ $\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$ will be partitioned into subsets in a corresponding way, such that roots in different subsets are mutually orthogonal.

Now the set of graphs which can occur as Dynkin diagrams of semisimple Lie algebras turns out to be quite restricted. In order to determine the possible Dynkin diagrams it is useful to introduce a quadratic form $\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)$ which is defined in terms of the Dynkin diagram. We define

$$
\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)=2 \sum_{i=1}^{l} x_{i}^{2}-\sum_{\substack{i, j=1 \\ i \neq j}}^{l} \sqrt{ } n_{i j} x_{i} x_{j} .
$$

We illustrate this definition in the cases $l=1,2$.

## Dynkin diagram Quadratic form

| $\circ$ | $2 x_{1}^{2}$ |
| :---: | :---: |
| $\circ$ | $\circ$ |$c$| $\circ x_{1}^{2}+2 x_{2}^{2}$ |  |
| :---: | :---: |
| $\circ$ | $\circ$ | | $2 x_{1}^{2}-2 x_{1} x_{2}+2 x_{2}^{2}$ |
| :---: |
| $\rightleftharpoons$ |

Proposition 6.6 The quadratic form $\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)$ is positive definite.

Proof. We have, for $i \neq j$,

$$
n_{i j}=A_{i j} A_{j i}=2 \frac{\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} \cdot 2 \frac{\left\langle\alpha_{j}, \alpha_{i}\right\rangle}{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}
$$

hence $-\sqrt{ } n_{i j}=2 \frac{\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\left|\alpha_{i}\right|\left|\alpha_{j}\right|}$ since $\left\langle\alpha_{i}, \alpha_{j}\right\rangle \leq 0$. For $i=j$ we have $\frac{2\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\left|\alpha_{i}\right|\left|\alpha_{j}\right|}=2$. Thus the quadratic form may be written

$$
\begin{aligned}
\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right) & =\sum_{i, j=1}^{l} \frac{2\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\left|\alpha_{i}\right|\left|\alpha_{j}\right|} x_{i} x_{j}=2\left\langle\sum_{i=1}^{l} \frac{x_{i} \alpha_{i}}{\left|\alpha_{i}\right|}, \sum_{j=1}^{l} \frac{x_{j} \alpha_{j}}{\left|\alpha_{j}\right|}\right\rangle \\
& =2\langle y, y\rangle \quad \text { where } y=\sum_{i=1}^{l} \frac{x_{i} \alpha_{i}}{\left|\alpha_{i}\right|}
\end{aligned}
$$

Thus $\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right) \geq 0$. Moreover if $\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)=0$ then $y=0$. Since $\alpha_{1}, \ldots, \alpha_{l}$ are linearly independent this implies that $x_{i}=0$ for all $i$. Thus the quadratic form is positive definite.

Now the connected components of the Dynkin diagram of any semisimple Lie algebra satisfy the following conditions:
(A) The graph is connected.
(B) Any pair of distinct vertices are joined by $0,1,2$ or 3 edges.
(C) The corresponding quadratic form $\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)$ is positive definite.

We shall approach the problem of finding the possible Dynkin diagrams by determining all graphs satisfying conditions (A), (B), (C). Having determined all such graphs we shall consider subsequently which ones occur as Dynkin diagrams.

### 6.3 Classification of Dynkin diagrams

The main result which we shall obtain in this section is as follows.

Theorem 6.7 The graphs satisfying conditions (A), (B), (C) shown in Section 6.2 are just those in the following list.


Proof. We shall show first that the graphs on this list satisfy conditions (A), (B), (C). It is obvious that they satisfy (A) and (B). We shall therefore concentrate on condition (C).

We recall from linear algebra that a quadratic form $\sum a_{i j} x_{i} x_{j}$ is positive definite if and only if all the leading minors of its symmetric matrix $\left(a_{i j}\right)$ have positive determinant. This condition is

$$
\left|a_{11}\right|>0,\left|\begin{array}{ll}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{array}\right|>0, \ldots, \operatorname{det}\left(a_{i j}\right)>0
$$

Given a graph $\Gamma$ with $l$ vertices on the list in Theorem 6.7 we shall show that $\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)$ is positive definite by induction on $l$. If $l=1$ then $\Gamma=A_{1}$ and $\mathrm{Q}\left(x_{1}\right)=2 x_{1}^{2}$ is positive definite. If $l=2$ then $\Gamma$ is $A_{2}, B_{2}$ or $G_{2}$. The symmetric matrix representing $\mathrm{Q}\left(x_{1}, x_{2}\right)$ is then

$$
\begin{array}{ccc}
\left(\begin{array}{cc}
2 & -1 \\
-1 & 2
\end{array}\right) & \left(\begin{array}{cc}
2 & -\sqrt{ } 2 \\
-\sqrt{ } 2 & 2
\end{array}\right) & \left(\begin{array}{cc}
2 & -\sqrt{ } 3 \\
-\sqrt{ } 3 & 2
\end{array}\right) \\
A_{2} & B_{2} & G_{2}
\end{array}
$$

In these cases the leading minors have positive determinant.
Now assume $l \geq 3$. Then inspection of the list of graphs in Theorem 6.7 shows that $\Gamma$ contains at least one vertex which is joined to just one other vertex of $\Gamma$, and joined to it by a single edge. Let such a vertex be labelled $l$, and let the vertex it is joined to be labelled $l-1$. We write $\Gamma=\Gamma_{l}$, and the graph obtained from $\Gamma_{l}$ by removing the vertex $l$ by $\Gamma_{l-1}$, and the graph obtained from $\Gamma_{l-1}$ by removing the vertex $l-1$ by $\Gamma_{l-2}$. Let $\operatorname{det} \Gamma_{l}$ be the determinant of the symmetric matrix representing the quadratic form $\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)$ associated
to $\Gamma_{l}$. We observe from the list of graphs that $\Gamma_{l-1}$ and $\Gamma_{l-2}$ also lie in the list. Moreover we have

$$
\operatorname{det} \Gamma_{l}=\left|\begin{array}{cccc} 
& & & \\
& & & \\
& & & \\
& & & 0 \\
& & & \\
& & & -1 \\
0 & \ldots & 0 & -1
\end{array}\right|=2 \operatorname{det} \Gamma_{l-1}-\operatorname{det} \Gamma_{l-2}
$$

by expanding the determinant by its last row. This gives us an inductive way of calculating $\operatorname{det} \Gamma_{l}$. In particular we have

$$
\operatorname{det} A_{1}=2, \quad \operatorname{det} A_{2}=3, \quad \operatorname{det} A_{l}=2 \operatorname{det} A_{l-1}-\operatorname{det} A_{l-2}
$$

Thus $\operatorname{det} A_{l}=l+1$.

$$
\operatorname{det} A_{1}=2, \quad \operatorname{det} B_{2}=2, \quad \operatorname{det} B_{3}=2, \quad \operatorname{det} B_{l}=2 \operatorname{det} B_{l-1}-\operatorname{det} B_{l-2}
$$

Thus $\operatorname{det} B_{l}=2$.

$$
\operatorname{det} A_{3}=4, \quad \operatorname{det} D_{4}=4, \quad \operatorname{det} D_{5}=4, \quad \operatorname{det} D_{l}=2 \operatorname{det} D_{l-1}-\operatorname{det} D_{l-2}
$$

Thus $\operatorname{det} D_{l}=4$.

$$
\begin{aligned}
& \operatorname{det} E_{6}=2 \operatorname{det} D_{5}-\operatorname{det} A_{4}=3 \\
& \operatorname{det} E_{7}=2 \operatorname{det} D_{6}-\operatorname{det} A_{5}=2 \\
& \operatorname{det} E_{8}=2 \operatorname{det} D_{7}-\operatorname{det} A_{6}=1 \\
& \operatorname{det} F_{4}=2 \operatorname{det} B_{3}-\operatorname{det} A_{2}=1
\end{aligned}
$$

Thus we have shown that $\operatorname{det} \Gamma_{l}>0$ for all $\Gamma_{l}$
Now the leading minors of the symmetric matrix associated to $\Gamma_{l}$ are the symmetric matrices associated to certain subgraphs of $\Gamma_{l}$. The numbering can be chosen so that all these subgraphs are connected. However, the list of graphs has the property that any connected subgraph of a graph on the list is also on the list. Thus the determinant of every leading minor of the given symmetric matrix is positive. Hence the quadratic form $\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)$ associated to $\Gamma_{l}$ is positive definite.

Thus we have shown that the graphs on our list satisfy conditions (A), (B), (C). We wish to prove the converse, i.e. that any graph satisfying conditions (A), (B), (C) is on our list. Before being able to prove this we shall need some lemmas.

Lemma 6.8 For each of the graphs on the following list the corresponding quadratic form $\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)$ has determinant 0 .

$\tilde{A}_{2}$

$\tilde{A}_{3}$





$\tilde{D}_{4}$



$\circ-\infty$



Proof. First consider the graphs $\Gamma=\tilde{A}_{l}$. Each row of the symmetric matrix of the given quadratic form contains one entry 2 , two entries -1 , and remaining entries 0 . Thus the sum of the columns is zero and $\operatorname{det} \tilde{A}_{l}=0$.

In all the other graphs $\Gamma$ on the list we can find a vertex $l$ joined to just one other vertex $l-1$. Moreover $l$ is joined to $l-1$ by a single edge or a double edge. If there is a single edge we may use the formula

$$
\operatorname{det} \Gamma_{l}=2 \operatorname{det} \Gamma_{l-1}-\operatorname{det} \Gamma_{l-2}
$$

as before. If there is a double edge we obtain instead

$$
\operatorname{det} \Gamma_{l}=2 \operatorname{det} \Gamma_{l-1}-2 \operatorname{det} \Gamma_{l-2}
$$

We may use these formulae to calculate all the determinants inductively.

$$
\begin{array}{ll}
\operatorname{det} \tilde{B}_{3}=2 \operatorname{det} A_{3}-2\left(\operatorname{det} A_{1}\right)^{2}=0 \\
\operatorname{det} \tilde{B}_{l}=2 \operatorname{det} D_{l}-2 \operatorname{det} D_{l-1}=0 & \text { for } l \geq 4 \\
\operatorname{det} \tilde{C}_{2}=2 \operatorname{det} B_{2}-2 \operatorname{det} A_{1}=0 & \\
\operatorname{det} \tilde{C}_{l}=2 \operatorname{det} B_{l}-2 \operatorname{det} B_{l-1}=0 & \text { for } l \geq 3 \\
\operatorname{det} \tilde{D}_{4}=2 \operatorname{det} D_{4}-\left(\operatorname{det} A_{1}\right)^{3}=0 &
\end{array}
$$

$$
\begin{aligned}
& \operatorname{det} \tilde{D}_{l}=2 \operatorname{det} D_{l}-\operatorname{det} D_{l-2} \cdot \operatorname{det} A_{1}=0 \quad \text { for } l \geq 5 \\
& \operatorname{det} \tilde{E}_{6}=2 \operatorname{det} E_{6}-\operatorname{det} A_{5}=0 \\
& \operatorname{det} \tilde{E}_{7}=2 \operatorname{det} E_{7}-\operatorname{det} D_{6}=0 \\
& \operatorname{det} \tilde{E}_{8}=2 \operatorname{det} E_{8}-\operatorname{det} E_{7}=0 \\
& \operatorname{det} \tilde{F}_{4}=2 \operatorname{det} F_{4}-\operatorname{det} B_{3}=0 \\
& \operatorname{det} \tilde{G}_{2}=2 \operatorname{det} G_{2}-\operatorname{det} A_{1}=0
\end{aligned}
$$

Lemma 6.9 Let $\Gamma$ be a graph satisfying conditions $(A),(B),(C)$ and $\Gamma^{\prime}$ be a connected graph obtained from $\Gamma$ by omitting vertices or decreasing the number of edges between vertices or both. Then $\Gamma^{\prime}$ satisfies conditions (A), (B), (C) also.

Proof. $\Gamma^{\prime}$ clearly satisfies (A) and (B). We must show it satisfies (C). Let $\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)$ be the quadratic form of $\Gamma$ and $\mathrm{Q}^{\prime}\left(x_{1}, \ldots, x_{m}\right)$ be the quadratic form of $\Gamma^{\prime}$, where $m \leq l$. We have

$$
\begin{aligned}
& \mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)=2 \sum_{i=1}^{l} x_{i}^{2}-\sum_{\substack{i, j=1 \\
i \neq j}}^{l} \sqrt{ } n_{i j} x_{i} x_{j} \\
& \mathrm{Q}^{\prime}\left(x_{1}, \ldots, x_{m}\right)=2 \sum_{i=1}^{m} x_{i}^{2}-\sum_{\substack{i, j=1 \\
i \neq j}}^{m} \sqrt{ } n_{i j}^{\prime} x_{i} x_{j}
\end{aligned}
$$

where $n_{i j}^{\prime} \leq n_{i j}$ for $i, j \in\{1, \ldots, m\}$. Suppose if possible that $\mathrm{Q}^{\prime}$ is not positive definite. Then there exist $y_{1}, \ldots, y_{m} \in \mathbb{R}$, not all zero, with $\mathrm{Q}^{\prime}\left(y_{1}, \ldots, y_{m}\right) \leq 0$. Consider $\mathrm{Q}\left(\left|y_{1}\right|, \ldots,\left|y_{m}\right|, 0, \ldots, 0\right)$. We have

$$
\begin{aligned}
\mathrm{Q}\left(\left|y_{1}\right|, \ldots,\left|y_{m}\right|, 0, \ldots, 0\right) & =2 \sum_{i=1}^{m}\left|y_{i}\right|^{2}-\sum_{\substack{i, j=1 \\
i \neq j}}^{m} \sqrt{ } n_{i j}\left|y_{i}\right|\left|y_{j}\right| \\
& \leq 2 \sum_{i=1}^{m} y_{i}^{2}-\sum_{\substack{i, j=1 \\
i \neq j}}^{m} \sqrt{ } n_{i j}^{\prime}\left|y_{i}\right|\left|y_{j}\right| \\
& \leq 2 \sum_{i=1}^{m} y_{i}^{2}-\sum_{\substack{i, j=1 \\
i \neq j}}^{m} \sqrt{ } n_{i j}^{\prime} y_{i} y_{j} \\
& =\mathrm{Q}^{\prime}\left(y_{1}, \ldots, y_{m}\right) \leq 0 .
\end{aligned}
$$

Hence $\mathrm{Q}\left(\left|y_{1}\right|, \ldots,\left|y_{m}\right|, 0, \ldots, 0\right) \leq 0$ but $\left(\left|y_{1}\right|, \ldots,\left|y_{m}\right|, 0, \ldots, 0\right)$ is not the zero vector. This contradicts the fact that $\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)$ is positive definite. Hence $\mathrm{Q}^{\prime}\left(x_{1}, \ldots, x_{m}\right)$ must be positive definite also.

Having Lemmas 6.8 and 6.9 at our disposal we are now able to complete the proof of Theorem 6.7.

Let $\Gamma$ be a graph satisfying conditions (A), (B), (C). Then, by Lemmas 6.8 and $6.9, \Gamma$ can have no subgraph of type $\tilde{A}_{l}, \tilde{B}_{l}, \tilde{C}_{l}, \tilde{D}_{l}, \tilde{E}_{6}, \tilde{E}_{7}, \tilde{E}_{8}, \tilde{F}_{4}$ or $\tilde{G}_{2}$. (By a subgraph of $\Gamma$ we mean a graph obtainable from $\Gamma$ by removing vertices, or removing edges, or both.) We shall use this information to show that $\Gamma$ must be one of the graphs on the list in Theorem 6.7.

In the first place we see that $\Gamma$ contains no cycles, otherwise it would contain a subgraph of type $\tilde{A}_{l}$ for some $l \geq 2$.

Suppose that $\Gamma$ contains a triple edge. Then $\Gamma$ must be the graph $G_{2}$, otherwise $\Gamma$ would contain a subgraph $\tilde{G}_{2}$.
Thus we may assume that $\Gamma$ contains no triple edge. Suppose $\Gamma$ contains a double edge. Then $\Gamma$ cannot have more than one double edge, otherwise it would contain a subgraph $\tilde{C}_{l}$ for some $l \geq 2$. Now $\Gamma$ cannot contain a branch point in addition to a double edge, as otherwise it would contain a subgraph $\tilde{B}_{l}$ for some $l \geq 3$. Thus $\Gamma$ is a chain containing just one double edge. If the double edge occurs at one end of the chain then $\Gamma=B_{l}$ for some $l \geq 2$. If not then we must have $\Gamma=F_{4}$, since otherwise $\Gamma$ would contain a subgraph $\tilde{F}_{4}$.

Thus we may assume that $\Gamma$ contains no double or triple edges. If $\Gamma$ contains no branch point then $\Gamma=A_{l}$ for some $l \geq 1$. Thus we suppose that $\Gamma$ contains at least one branch point. Now $\Gamma$ cannot contain more than one branch point, as otherwise it would contain a subgraph $\tilde{D}_{l}$ for some $l \geq 5$. Thus $\Gamma$ contains exactly one branch point. There must be exactly three branches emerging from this branch point, since otherwise $\Gamma$ would contain a subgraph $\tilde{D}_{4}$. Let the number of vertices on the three branches be $l_{1}, l_{2}, l_{3}$ with $l_{1} \geq l_{2} \geq l_{3}$. Then the total number of vertices of $\Gamma$ is $l=l_{1}+l_{2}+l_{3}+1$.
Now we must have $l_{3}=1$, as otherwise we have $l_{i} \geq 2$ for $i=1,2,3$ and $\Gamma$ contains a subgraph $\tilde{E}_{6}$. If $l_{2}=1$ then $\Gamma=D_{l}$ for some $l \geq 4$. Thus we may assume $l_{2} \geq 2$. In fact we must have $l_{2}=2$, as otherwise we have $l_{1} \geq 3, l_{2} \geq 3$ and $\Gamma$ contains a subgraph $\tilde{E}_{7}$. Thus we may assume $l_{3}=1, l_{2}=2$. We must have $l_{1} \leq 4$ since otherwise $\Gamma$ contains a subgraph $\tilde{E}_{8}$. Thus $\Gamma$ has type $E_{6}, E_{7}$ or $E_{8}$.

Thus we have now determined all possibilities for $\Gamma$, and seen that $\Gamma$ must be one of the graphs which appear on the list in Theorem 6.7. This completes the proof.

Corollary 6.10 Let $\Delta$ be the Dynkin diagram of a semisimple Lie algebra. Then each connected component of $\Delta$ must be one of the graphs

$$
A_{l}, \quad l \geq 1 ; B_{l}, \quad l \geq 2 ; D_{l}, \quad l \geq 4 ; E_{6} ; E_{7} ; E_{8} ; F_{4} ; G_{2}
$$

We shall consider later whether all these graphs actually occur as Dynkin diagrams.

### 6.4 Classification of Cartan matrices

We recall that the Dynkin diagram is determined by the Cartan matrix by the property

$$
n_{i j}=A_{i j} A_{j i} \quad i \neq j
$$

However, the Cartan matrix is not always uniquely determined by the Dynkin diagram. If we know the integers $n_{i j} \in\{0,1,2,3\}$ for all $i, j$ with $i \neq j$ we consider to what extent the $A_{i j}$ are determined. If $n_{i j}=0$ then we must have $A_{i j}=0$ and $A_{j i}=0$ since $A_{i j}=0$ if and only if $A_{j i}=0$. If $n_{i j}=1$ then we must have $A_{i j}=-1$ and $A_{j i}=-1$ since $A_{i j} \in \mathbb{Z}, A_{j i} \in \mathbb{Z}, A_{i j} \leq 0, A_{j i} \leq 0$. However, if $n_{i j}=2$ there are two possibilities for the factorisation $n_{i j}=A_{i j} A_{j i}$. Either we have $2=-1 \cdot-2$ or $2=-2 \cdot-1$. Thus we have either $A_{i j}=-1, A_{j i}=-2$ or $A_{i j}=-2, A_{j i}=-1$. Similarly if $n_{i j}=3$ we have either $A_{i j}=-1, A_{j i}=-3$ or $A_{i j}=-3, A_{j i}=-1$.

In the connected graphs in Corollary 6.10 the only ones which give rise to such an ambiguity are $B_{l}, \quad l \geq 2 ; F_{4}$ and $G_{2}$. In these graphs we shall place an arrow on the double or triple edges. The direction of the arrow is determined as follows. The arrow points from vertex $i$ to vertex $j$ if and only if $\left|\alpha_{i}\right|>\left|\alpha_{j}\right|$, that is $\left|A_{j i}\right|>\left|A_{i j}\right|$.

Thus in the situation

we have $\left|\alpha_{i}\right|=\sqrt{ } 2\left|\alpha_{j}\right|, \quad A_{i j}=-1, \quad A_{j i}=-2$. In the situation

we have $\left|\alpha_{i}\right|=\sqrt{ } 3\left|\alpha_{j}\right|, \quad A_{i j}=-1, \quad A_{j i}=-3$. The arrow may thus be regarded as an inequality sign on the lengths of the fundamental roots at the vertices.

The set of possible connected Dynkin diagrams, including arrows, is shown on the following standard list.

### 6.11 Standard list of connected Dynkin diagrams



We note that, since the diagrams of types $\mathrm{B}_{2}, \mathrm{~F}_{4}, \mathrm{G}_{2}$ are symmetric, it does not matter in which direction the arrow is drawn in these cases.

The connected components of the Dynkin diagram of any semisimple Lie algebra must appear on this standard list.

We next obtain a standard list of corresponding Cartan matrices. We say that two Cartan matrices $\left(A_{i j}\right),\left(A_{i j}^{\prime}\right)$ are equivalent if they have the same degree $l$ and there is a permutation $\sigma$ of $1, \ldots, l$ such that

$$
A_{i j}^{\prime}=A_{\sigma(i) \sigma(j)} .
$$

Equivalent Cartan matrices come from different labellings of the same Dynkin diagram. For each Dynkin diagram on the standard list 6.11 we choose a labelling and obtain a corresponding Cartan matrix which is uniquely determined. These Cartan matrices appear on the following list.
6.12 Standard list of indecomposable Cartan matrices

$$
\begin{aligned}
& E_{6}=\left(\begin{array}{cccccc}
2 & -1 & & & & \\
-1 & 2 & -1 & & & \\
& -1 & 2 & -1 & -1 & \\
& & -1 & 2 & & \\
& & -1 & & 2 & -1 \\
& & & & -1 & 2
\end{array}\right) . \\
& E_{7}=\left(\begin{array}{ccccccc}
2 & -1 & & & & & \\
-1 & 2 & -1 & & & & \\
& -1 & 2 & -1 & & & \\
& & -1 & 2 & -1 & -1 & \\
& & & -1 & 2 & & \\
& & & & & & \\
& & & & & & \\
& & & & & & \\
& & & & & &
\end{array}\right) \text {. } \\
& E_{8}=\left(\begin{array}{ccccccccc}
2 & -1 & & & & & & & \\
-1 & 2 & -1 & & & & & & \\
& -1 & 2 & -1 & & & & & \\
& & -1 & 2 & -1 & & & & \\
& & & -1 & 2 & -1 & -1 & \\
& & & & -1 & 2 & & \\
& & & & & -1 & & 2 & -1 \\
& & & & & & & & -1
\end{array}\right) . \\
& F_{4}=\left(\begin{array}{cccc}
2 & -1 & & \\
-1 & 2 & -1 & \\
& -2 & 2 & -1 \\
& & -1 & 2
\end{array}\right) . \\
& G_{2}=\left(\begin{array}{cc}
2 & -1 \\
-3 & 2
\end{array}\right) .
\end{aligned}
$$

A Cartan matrix is called indecomposable if its Dynkin diagram is connected. Any Cartan matrix will determine a set of indecomposable Cartan matrices, unique up to equivalence, whose Dynkin diagrams are the connected components of the Dynkin diagram of the given Cartan matrix.

If $A$ is the Cartan matrix of any semisimple Lie algebra, each indecomposable component of $A$ will be equivalent to some Cartan matrix from the above standard list.

Proposition 6.13 If a semisimple Lie algebra L has a connected Dynkin diagram then $L$ is simple.

Proof. Let $L=H \oplus \sum_{\alpha \in \Phi} L_{\alpha}$ be a Cartan decomposition giving rise to the Dynkin diagram $\Delta$. Let $I$ be a non-zero ideal of $L$. We shall show that $I=L$, thus proving that $L$ is simple.

We first aim to prove that $I \cap H \neq O$. Suppose if possible that $I \cap H=O$. Let $e_{\alpha}$ be a non-zero element of $L_{\alpha}$ and choose a non-zero element $x \in I$ with

$$
x=h+\sum_{\alpha \in \Phi} \mu_{\alpha} e_{\alpha} \quad h \in H, \mu_{\alpha} \in \mathbb{C}
$$

such that the number of non-zero $\mu_{\alpha}$ is as small as possible. Since $I \cap H=O$ there exists some $\mu_{\beta} \neq 0$. Then we have

$$
\left[h_{\beta}^{\prime} x\right]=\sum_{\alpha \in \Phi} \mu_{\alpha}\left[h_{\beta}^{\prime} e_{\alpha}\right]=\sum_{\alpha \in \Phi} \mu_{\alpha} \alpha\left(h_{\beta}^{\prime}\right) e_{\alpha}
$$

Now by Proposition 4.18 we can choose $e_{\beta} \in L_{\beta}$ and $e_{-\beta} \in L_{-\beta}$ such that $\left[e_{\beta} e_{-\beta}\right]=h_{\beta}^{\prime}$. Thus $\left[\left[h_{\beta}^{\prime} x\right] e_{-\beta}\right]=\sum_{\alpha \in \Phi} \mu_{\alpha} \alpha\left(h_{\beta}^{\prime}\right)\left[e_{\alpha} e_{-\beta}\right]=\mu_{\beta} \beta\left(h_{\beta}^{\prime}\right) h_{\beta}^{\prime}+$ $\sum_{\substack{\alpha \in \Phi \\ \alpha \neq \beta}} \mu_{\alpha} \alpha\left(h_{\beta}^{\prime}\right) N_{\alpha,-\beta} e_{\alpha-\beta}$ where $\left[e_{\alpha} e_{-\beta}\right]=N_{\alpha,-\beta} e_{\alpha-\beta}$. Now we have $\left[\left[h_{\beta}^{\prime} x\right] e_{-\beta}\right] \in I$ since $x \in I$ and $\left[\left[h_{\beta}^{\prime} x\right] e_{-\beta}\right] \neq 0$ since $\mu_{\beta} \neq 0$ and $\beta\left(h_{\beta}^{\prime}\right)=$ $\left\langle h_{\beta}^{\prime}, h_{\beta}^{\prime}\right\rangle \neq 0$. Moreover the number of non-zero terms coming from the root spaces $L_{\alpha}$ is less for $\left[\left[h_{\beta}^{\prime} x\right] e_{-\beta}\right]$ than it was for $x$. This contradicts the choice of $x$. We can therefore deduce that $I \cap H \neq O$.

The next step is to show that $I \supset H$. Suppose if possible this is not so. Then

$$
O \neq I \cap H \neq H
$$

This implies that there exist $\alpha_{i} \in \Pi$ and $x \in I \cap H$ such that $\left\langle h_{\alpha_{i}}^{\prime}, x\right\rangle \neq 0$. For if $I \cap H$ were orthogonal to each $h_{\alpha_{i}}^{\prime}$ it would be orthogonal to the whole of $H$ and would therefore be $O$. Then we have

$$
\left[x e_{\alpha_{i}}\right]=\alpha_{i}(x) e_{\alpha_{i}}=\left\langle h_{\alpha_{i}}^{\prime}, x\right\rangle e_{\alpha_{i}} \in I
$$

Since $\left\langle h_{\alpha_{i}}^{\prime}, x\right\rangle \neq 0$ we deduce that $e_{\alpha_{i}} \in I$. Thus $\left[e_{\alpha_{i}} e_{-\alpha_{i}}\right]=h_{\alpha_{i}}^{\prime} \in I$.
We can therefore divide the $\alpha_{i} \in \Pi$ into two classes, those with $h_{\alpha_{i}}^{\prime} \in I$ and those with $h_{\alpha_{i}}^{\prime} \notin I$. Both classes are non-empty. Furthermore if $h_{\alpha_{i}}^{\prime} \in I$ and $h_{\alpha_{j}}^{\prime} \notin I$ then $\left\langle h_{\alpha_{j}}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle=0$. This means that vertices $i, j$ are not joined in the Dynkin diagram $\Delta$, so $\Delta$ is disconnected. This is a contradiction, thus we deduce that $I \supset H$.

Finally we show that $I=L$. Let $\alpha \in \Phi$. Then we have

$$
\left[h_{\alpha}^{\prime} e_{\alpha}\right]=\alpha\left(h_{\alpha}^{\prime}\right) e_{\alpha}=\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle e_{\alpha}
$$

Since $h_{\alpha}^{\prime} \in I$ we have $\left[h_{\alpha}^{\prime} e_{\alpha}\right] \in I$, and since $\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle \neq 0$ we deduce that $e_{\alpha} \in I$. This is true for all $\alpha \in \Phi$ and so $I=L$. Thus $L$ is simple.

We next consider what happens when the Dynkin diagram of $L$ is disconnected.

We first define an action of the Weyl group on $H$. The Weyl group was introduced in Section 5.2 as a group of non-singular linear transformations on the real vector space $H_{\mathbb{R}}^{*}$. This action can be extended by linearity to give an action of $W$ on $H^{*}$ by $\mathbb{C}$-linear transformations. We also define an action of $W$ on $H$ by $h \rightarrow w h$ where

$$
\lambda(w h)=\left(w^{-1} \lambda\right) h \quad \text { for all } h \in H, \lambda \in H^{*}, w \in W
$$

There is a unique element $w h \in H$ satisfying this condition, and

$$
w_{1}\left(w_{2} h\right)=\left(w_{1} w_{2}\right) h \quad \text { for all } w_{1}, w_{2} \in W
$$

The actions of $W$ on $H^{*}$ and $H$ are compatible with the isomorphism $H^{*} \rightarrow H$ given by $\lambda \rightarrow h_{\lambda}^{\prime}$ where $\lambda(x)=\left\langle h_{\lambda}^{\prime}, x\right\rangle$ for all $x \in H$. For suppose $w(\lambda)=\mu$ for $\lambda, \mu \in H^{*}$. Then

$$
\begin{aligned}
\left\langle w\left(h_{\lambda}^{\prime}\right), x\right\rangle & =\left\langle h_{\lambda}^{\prime}, w^{-1}(x)\right\rangle=\lambda\left(w^{-1}(x)\right)=(w \lambda) x \\
& =\mu(x)=\left\langle h_{\mu}^{\prime}, x\right\rangle \quad \text { for all } x \in H .
\end{aligned}
$$

Hence $w(\lambda)=\mu$ implies $w\left(h_{\lambda}^{\prime}\right)=h_{\mu}^{\prime}$.
Since we know that

$$
s_{\alpha}(\lambda)=\lambda-2 \frac{\langle\alpha, \lambda\rangle}{\langle\alpha, \alpha\rangle} \alpha \quad \text { for } \alpha \in \Phi, \lambda \in H^{*}
$$

it follows that

$$
s_{\alpha}(x)=x-2 \frac{\left\langle h_{\alpha}^{\prime}, x\right\rangle}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle} h_{\alpha}^{\prime} \quad \text { for } x \in H
$$

Proposition 6.14 Let L be a semisimple Lie algebra whose Dynkin diagram $\Delta$ splits into connected components $\Delta_{1}, \ldots, \Delta_{r}$. Then we have

$$
L=L_{1} \oplus \cdots \oplus L_{r}
$$

a direct sum of Lie algebras, where $L_{i}$ is a simple Lie algebra with Dynkin diagram $\Delta_{i}$.

Proof. We have $\Delta=\Delta_{1} \dot{\cup} \Delta_{2} \dot{\cup} \cdots \dot{\cup} \Delta_{r}$. Let $\Pi_{i}$ be the subset of $\Pi$ corresponding to the vertices in $\Delta_{i}$. Then we have

$$
\Pi=\Pi_{1} \dot{\cup} \Pi_{2} \dot{\cup} \cdots \dot{\cup} \Pi_{r} .
$$

Moreover we have $\langle\alpha, \beta\rangle=0$ if $\alpha \in \Pi_{i}, \beta \in \Pi_{j}$ and $i \neq j$. Let $H_{i}$ be the subspace of $H$ spanned by the elements $h_{\alpha}^{\prime}$ with $\alpha \in \Pi_{i}$. Then we have

$$
H=H_{1} \oplus H_{2} \oplus \cdots \oplus H_{r}
$$

where $\left\langle h, h^{\prime}\right\rangle=0$ if $h \in H_{i}, h^{\prime} \in H_{j}$ and $i \neq j$.
Now let $\alpha \in \Pi_{i}$ and consider the fundamental reflection $s_{\alpha} \in W$. It is clear that $s_{\alpha}$ transforms $H_{i}$ into itself and fixes each vector in $H_{j}$ for all $j \neq i$. Thus we have

$$
s_{\alpha}\left(H_{j}\right)=H_{j} \quad j=1, \ldots, r
$$

Since the elements $s_{\alpha}$ generate the Weyl group $W$ we deduce that

$$
w\left(H_{j}\right)=H_{j} \quad j=1, \ldots, r \quad w \in W
$$

Now for all $\alpha \in \Phi$ we have $h_{\alpha}^{\prime}=w\left(h_{\alpha_{i}}^{\prime}\right)$ for some $\alpha_{i} \in \Pi$ and some $w \in W$, by Proposition 5.12 and the definition of the $W$-action on $H$. It follows that each $h_{\alpha}^{\prime}, \alpha \in \Phi$, lies in $H_{i}$ for some $i$. Let $\Phi_{i}$ be the set of all $\alpha \in \Phi$ such that $h_{\alpha}^{\prime} \in H_{i}$. Then we have

$$
\Phi=\Phi_{1} \dot{\cup} \Phi_{2} \dot{\cup} \cdots \dot{\cup} \Phi_{r}
$$

We define $L_{i}$ to be the subspace of $L$ spanned by $H_{i}$ and the $e_{\alpha}$ for all $\alpha \in \Phi_{i}$. We deduce from the Cartan decomposition of $L$ that

$$
L=L_{1} \oplus L_{2} \oplus \cdots \oplus L_{r}
$$

a direct sum of vector spaces. In fact we can see that each $L_{i}$ is a subalgebra of $L$. It is sufficient to verify that $\left[e_{\alpha} e_{\beta}\right] \in L_{i}$ if $\alpha, \beta \in \Phi_{i}$. If $\alpha+\beta \in \Phi$ then we have $\alpha+\beta \in \Phi_{i}$ since $h_{\alpha+\beta}^{\prime}=h_{\alpha}^{\prime}+h_{\beta}^{\prime} \in H_{i}$. If $\alpha+\beta=0$ then $\left[e_{\alpha} e_{\beta}\right]$ is a multiple of $h_{\alpha}^{\prime}$ and so lies in $H_{i}$, thus in $L_{i}$. If $\alpha+\beta$ is non-zero but not a root then $\left[e_{\alpha} e_{\beta}\right]=0$. In either case we have $\left[e_{\alpha} e_{\beta}\right] \in L_{i}$. Thus $L_{i}$ is a subalgebra of $L$.

We show next that $\left[L_{i} L_{j}\right]=O$ if $i \neq j$. Let $\alpha \in \Phi_{i}$ and $\beta \in \Phi_{j}$. Then we have

$$
\left[h_{\alpha}^{\prime} e_{\beta}\right]=\beta\left(h_{\alpha}^{\prime}\right) e_{\beta}=\left\langle h_{\beta}^{\prime}, h_{\alpha}^{\prime}\right\rangle e_{\beta}=0
$$

and similarly $\left[e_{\alpha} h_{\beta}^{\prime}\right]=0$. We also have $\left[e_{\alpha} e_{\beta}\right]=0$. For $\alpha+\beta \notin \Phi$ since $h_{\alpha}^{\prime}+h_{\beta}^{\prime}$ does not lie in any subspace $H_{k}$ of $H$. It follows that $\left[L_{i} L_{j}\right]=O$.

We now know that each $L_{i}$ is an ideal of $L$, since

$$
\left[L_{i} L\right] \subset \sum_{j}\left[L_{i} L_{j}\right] \subset\left[L_{i} L_{i}\right] \subset L_{i}
$$

This implies that

$$
\left[x_{1}+\cdots+x_{r}, y_{1}+\cdots+y_{r}\right]=\left[x_{1} y_{1}\right]+\cdots+\left[x_{r} y_{r}\right]
$$

where $x_{i}, y_{i} \in L_{i}$. Hence

$$
L=L_{1} \oplus L_{2} \oplus \cdots \oplus L_{r}
$$

is a direct sum of Lie algebras.
Now each $L_{i}$ is a semisimple Lie algebra. For let $I$ be a soluble ideal of $L_{i}$. Since $\left[I L_{j}\right]=O$ for all $j \neq i$ the ideal $I$ is an ideal of $L$. Since $L$ is semisimple we have $I=O$. Hence $L_{i}$ is semisimple.
We next observe that $H_{i}$ is a Cartan subalgebra of $L_{i}$. The subalgebra $H_{i}$ is abelian, hence nilpotent. Let $x \in L_{i}$ satisfy $x \in N\left(H_{i}\right)$. Then $[x h] \in H_{i}$ for all $h \in H_{i}$. We also have $[x h]=0$ for all $h \in H_{j}$ with $j \neq i$. It follows that $[x h] \in H$ for all $h \in H$. Since $H$ is a Cartan subalgebra of $L$ we have $N(H)=H$. Hence $x \in H$. Thus $x \in H \cap L_{i}=H_{i}$. Thus $H_{i}$ is a Cartan subalgebra of $L_{i}$.

We now consider the Cartan decomposition

$$
L_{i}=H_{i} \oplus \sum_{\alpha \in \Phi_{i}} \mathbb{C} e_{\alpha}
$$

of $L_{i}$ with respect to $H_{i}$. We see that $\Phi_{i}$ is the root system of $L_{i}$ with respect to $H_{i}$, that $\Pi_{i}$ is a fundamental system of roots in $\Phi_{i}$, and that $\Delta_{i}$ is the Dynkin diagram of $L_{i}$. Now $\Delta_{i}$ is connected. Thus the Lie algebra $L_{i}$ must be simple, by Proposition 6.13. Thus we have obtained a decomposition of $L$ as a direct sum of simple Lie algebras $L_{i}$, whose Dynkin diagrams are the connected components $\Delta_{i}$ of $\Delta$.

Corollary 6.15 A semisimple Lie algebra L has a connected Dynkin diagram if and only if $L$ is simple.

Proof. This follows from Propositions 6.13 and 6.14

## 7

## The existence and uniqueness theorems

We have seen that each non-trivial simple Lie algebra $L$ has a Dynkin diagram $\Delta$ which appears on the standard list 6.11 of connected Dynkin diagrams. In the present chapter we shall consider the converse question. Given a Dynkin diagram $\Delta$ on the standard list, is there a simple Lie algebra $L$ with Dynkin diagram $\Delta$ ? If so, is $L$ uniquely determined up to isomorphism? We shall show that both the existence and uniqueness properties hold. The proof of the uniqueness property is somewhat easier, and we shall prove this first. In order to do so we shall need some properties of the structure constants of the Lie algebra $L$.

### 7.1 Some properties of structure constants

Let $L$ be a simple Lie algebra with Dynkin diagram $\Delta$. Let $H$ be a Cartan subalgebra of $L$ and

$$
L=H \oplus \sum_{\alpha \in \Phi} L_{\alpha}
$$

be the Cartan decomposition of $L$ with respect to $H$. We know from Theorem 4.20 that $\operatorname{dim} L_{\alpha}=1$ for each $\alpha \in \Phi$. Let $e_{\alpha}$ be a non-zero element of $L_{\alpha}$. Let $\Pi$ be a fundamental system of roots in $\Phi$. Then the elements $h_{\alpha_{i}}^{\prime}$ for $\alpha_{i} \in \Pi$ form a basis for $H$. It will be convenient to choose a slightly different basis consisting of scalar multiples of the $h_{\alpha_{i}}^{\prime}$. We define $h_{i} \in H$ by

$$
h_{i}=\frac{2 h_{\alpha_{i}}^{\prime}}{\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle} .
$$

We note that $\alpha_{i}\left(h_{i}\right)=2$. Then $\left\{h_{i}, \quad i=1, \ldots, l ; e_{\alpha}, \quad \alpha \in \Phi\right\}$ is a basis of $L$. By Proposition 4.18 we know that $h_{\alpha}^{\prime} \in\left[L_{\alpha} L_{-\alpha}\right]$ for all $\alpha \in \Phi$. Thus,
if we have already chosen the $e_{\alpha}$ for all $\alpha \in \Phi^{+}$, we may choose the $e_{-\alpha}$ uniquely for $\alpha \in \Phi^{+}$to satisfy the condition

$$
\left[e_{\alpha} e_{-\alpha}\right]=\frac{2 h_{\alpha}^{\prime}}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle} .
$$

(This relation will then be satisfied for $\alpha \in \Phi^{-}$also.)
We define $h_{\alpha} \in H$ for each $\alpha \in \Phi$ by

$$
h_{\alpha}=\frac{2 h_{\alpha}^{\prime}}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle} .
$$

The element $h_{\alpha}$ is called the coroot corresponding to the root $\alpha$. In particular we have $h_{i}=h_{\alpha_{i}}$. We then have

$$
\left[e_{\alpha} e_{-\alpha}\right]=h_{\alpha} \quad \text { for all } \alpha \in \Phi
$$

We next consider the product $\left[e_{\alpha} e_{\beta}\right]$ when $\alpha+\beta \neq 0$. We have $\left[L_{\alpha} L_{\beta}\right]=O$ if $\alpha+\beta \neq 0$ and $\alpha+\beta \notin \Phi$. If $\alpha+\beta \in \Phi$ we have $\left[L_{\alpha} L_{\beta}\right] \subset L_{\alpha+\beta}$. We define $N_{\alpha, \beta} \in \mathbb{C}$ by the condition

$$
\left[e_{\alpha} e_{\beta}\right]=N_{\alpha, \beta} e_{\alpha+\beta}
$$

The numbers $N_{\alpha, \beta}$ for $\alpha, \beta, \alpha+\beta \in \Phi$ will be called the structure constants of $L$. They clearly depend upon the choice of the elements $e_{\alpha} \in L_{\alpha}$.

We now consider the multiplication of the basis vectors $\left\{h_{i}, e_{\alpha}\right\}$ of $L$. We have

$$
\begin{aligned}
& {\left[h_{i} h_{j}\right]=0} \\
& {\left[h_{i} e_{\alpha}\right]=\alpha\left(h_{i}\right) e_{\alpha}} \\
& {\left[e_{\alpha} e_{-\alpha}\right]=h_{\alpha}} \\
& {\left[e_{\alpha} e_{\beta}\right]=N_{\alpha, \beta} e_{\alpha+\beta} \quad \text { if } \alpha, \beta, \alpha+\beta \in \Phi} \\
& {\left[e_{\alpha} e_{\beta}\right]=0 \quad \text { if } \alpha+\beta \neq 0 \text { and } \alpha+\beta \notin \Phi .}
\end{aligned}
$$

In order to express $\left[e_{\alpha} e_{-\alpha}\right]$ as a linear combination of basis elements we may express $h_{\alpha}^{\prime}$ as a linear combination of the $h_{\alpha_{i}}^{\prime}, \alpha_{i} \in \Pi$, and so also express $h_{\alpha}$ as a linear combination of the $h_{i}$.

We shall now derive some relations between the structure constants $N_{\alpha, \beta}$.

Proposition 7.1 The structure constants $N_{\alpha, \beta}$ satisfy the following relations.
(i) $N_{\beta, \alpha}=-N_{\alpha, \beta}$.
(ii) If $\alpha, \beta, \gamma \in \Phi$ satisfy $\alpha+\beta+\gamma=0$ then $\frac{N_{\alpha, \beta}}{\langle\gamma, \gamma\rangle}=\frac{N_{\beta, \gamma}}{\langle\alpha, \alpha\rangle}=\frac{N_{\gamma, \alpha}}{\langle\beta, \beta\rangle}$.
(iii) $N_{\alpha, \beta} N_{-\alpha,-\beta}=-(p+1)^{2}$ where the $\alpha$-chain of roots through $\beta$ is $-p \alpha+\beta, \ldots, \beta, \ldots, q \alpha+\beta$.
(iv) If $\alpha, \beta, \gamma, \delta \in \Phi$ satisfy $\alpha+\beta+\gamma+\delta=0$ and no pair are negatives of one another, then

$$
\frac{N_{\alpha, \beta} N_{\gamma, \delta}}{\langle\alpha+\beta, \alpha+\beta\rangle}+\frac{N_{\beta, \gamma} N_{\alpha, \delta}}{\langle\beta+\gamma, \beta+\gamma\rangle}+\frac{N_{\gamma, \alpha} N_{\beta, \delta}}{\langle\gamma+\alpha, \gamma+\alpha\rangle}=0 .
$$

Proof. (i) This relation is clear.
(ii) Suppose $\alpha+\beta+\gamma=0$. We consider the Jacobi identity

$$
\left[\left[e_{\alpha} e_{\beta}\right] e_{\gamma}\right]+\left[\left[e_{\beta} e_{\gamma}\right] e_{\alpha}\right]+\left[\left[e_{\gamma} e_{\alpha}\right] e_{\beta}\right]=0
$$

This gives

$$
N_{\alpha, \beta}\left[e_{\alpha+\beta} e_{-(\alpha+\beta)}\right]+N_{\beta, \gamma}\left[e_{-\alpha} e_{\alpha}\right]+N_{\gamma, \alpha}\left[e_{-\beta} e_{\beta}\right]=0
$$

that is

$$
2 N_{\alpha, \beta} \frac{h_{\alpha+\beta}^{\prime}}{\left\langle h_{\alpha+\beta}^{\prime}, h_{\alpha+\beta}^{\prime}\right\rangle}=2 N_{\beta, \gamma} \frac{h_{\alpha}^{\prime}}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}+2 N_{\gamma, \alpha} \frac{h_{\beta}^{\prime}}{\left\langle h_{\beta}^{\prime}, h_{\beta}^{\prime}\right\rangle}
$$

Now the roots $\alpha, \beta$ are linearly independent since, if they were not, $\alpha+\beta$ could not be a root. Thus $h_{\alpha}^{\prime}, h_{\beta}^{\prime}$ are linearly independent and $h_{\alpha+\beta}^{\prime}=h_{\alpha}^{\prime}+h_{\beta}^{\prime}$. We deduce that

$$
\frac{N_{\alpha, \beta}}{\left\langle h_{\alpha+\beta}^{\prime}, h_{\alpha+\beta}^{\prime}\right\rangle}=\frac{N_{\beta, \gamma}}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}=\frac{N_{\gamma, \alpha}}{\left\langle h_{\beta}^{\prime}, h_{\beta}^{\prime}\right\rangle}
$$

that is

$$
\frac{N_{\alpha, \beta}}{\langle\gamma, \gamma\rangle}=\frac{N_{\beta, \gamma}}{\langle\alpha, \alpha\rangle}=\frac{N_{\gamma, \alpha}}{\langle\beta, \beta\rangle}
$$

(iii) Now suppose $\alpha, \beta \in \Phi$ are linearly independent. We consider the Jacobi identity

$$
\left[\left[e_{\alpha} e_{-\alpha}\right] e_{\beta}\right]+\left[\left[e_{-\alpha} e_{\beta}\right] e_{\alpha}\right]+\left[\left[e_{\beta} e_{\alpha}\right] e_{-\alpha}\right]=0
$$

This gives

$$
2 \frac{\left[h_{\alpha}^{\prime} e_{\beta}\right]}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}+N_{-\alpha, \beta} N_{-\alpha+\beta, \alpha} e_{\beta}+N_{\beta, \alpha} N_{\alpha+\beta,-\alpha} e_{\beta}=0
$$

We deduce that

$$
2 \frac{\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}+N_{-\alpha, \beta} N_{-\alpha+\beta, \alpha}+N_{\beta, \alpha} N_{\alpha+\beta,-\alpha}=0
$$

Using relations (i) and (ii) this may be written

$$
\begin{aligned}
& N_{\alpha, \beta} N_{-\alpha,-\beta} \frac{\langle\beta, \beta\rangle}{\langle\alpha+\beta, \alpha+\beta\rangle}-N_{\alpha,-\alpha+\beta} N_{-\alpha, \alpha-\beta} \frac{\langle-\alpha+\beta,-\alpha+\beta\rangle}{\langle\beta, \beta\rangle} \\
& \quad=2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle} .
\end{aligned}
$$

(If $-\alpha+\beta$ is not a root $N_{-\alpha, \beta}$ is interpreted as 0 so the middle term disappears.) We now consider the $\alpha$-chain of roots through $\beta$. Let it be

$$
-p \alpha+\beta, \ldots, \beta, \ldots, q \alpha+\beta
$$

We apply the same formula to the pairs $(\alpha, \beta)(\alpha,-\alpha+\beta) \ldots(\alpha,-p \alpha+\beta)$ and obtain

$$
\begin{aligned}
& N_{\alpha, \beta} N_{-\alpha,-\beta} \frac{\langle\beta, \beta\rangle}{\langle\alpha+\beta, \alpha+\beta\rangle}-N_{\alpha,-\alpha+\beta} N_{-\alpha, \alpha-\beta} \frac{\langle-\alpha+\beta,-\alpha+\beta\rangle}{\langle\beta, \beta\rangle}=2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle} \\
& N_{\alpha,-\alpha+\beta} N_{-\alpha, \alpha-\beta} \frac{\langle-\alpha+\beta,-\alpha+\beta\rangle}{\langle\beta, \beta\rangle}-N_{\alpha,-2 \alpha+\beta} N_{-\alpha, 2 \alpha-\beta} \frac{\langle-2 \alpha+\beta,-2 \alpha+\beta\rangle}{\langle-\alpha+\beta,-\alpha+\beta\rangle} \\
& \quad=2 \frac{\langle\alpha,-\alpha+\beta\rangle}{\langle\alpha, \alpha\rangle}
\end{aligned}
$$

$$
N_{\alpha,-p \alpha+\beta} N_{-\alpha, p \alpha-\beta} \frac{\langle-p \alpha+\beta,-p \alpha+\beta\rangle}{\langle-(p-1) \alpha+\beta,-(p-1) \alpha+\beta\rangle}=\frac{2\langle\alpha,-p \alpha+\beta\rangle}{\langle\alpha, \alpha\rangle} .
$$

(The last equation has only one term on the left since $-(p+1) \alpha+\beta$ is not a root.) Adding these equations we obtain

$$
N_{\alpha, \beta} N_{-\alpha,-\beta} \frac{\langle\beta, \beta\rangle}{\langle\alpha+\beta, \alpha+\beta\rangle}=2(p+1) \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle}-2 \frac{p(p+1)}{2} .
$$

However, we know from Proposition 4.22 that $2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle}=p-q$. Thus we have

$$
N_{\alpha, \beta} N_{-\alpha,-\beta} \frac{\langle\beta, \beta\rangle}{\langle\alpha+\beta, \alpha+\beta\rangle}=-(p+1) q .
$$

In order to obtain the required result $N_{\alpha, \beta} N_{-\alpha,-\beta}=-(p+1)^{2}$ we must show

$$
\frac{\langle\alpha+\beta, \alpha+\beta\rangle}{\langle\beta, \beta\rangle}=\frac{p+1}{q} .
$$

We recall from the proof of Proposition 6.1 that

$$
2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle} \cdot 2 \frac{\langle\beta, \alpha\rangle}{\langle\beta, \beta\rangle}=4 \cos ^{2} \theta
$$

where $\theta$ is the angle between $\alpha, \beta$ and hence that $2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle} \in\{0,-1,-2,-3\}$. Also from Proposition 4.22 we know that $2 \frac{\langle\alpha, \beta\rangle}{\langle\alpha, \alpha\rangle}=p-q$. If we choose $\beta$ to be the initial root in its $\alpha$-chain we have $p=0$ and hence $q \leq 3$. This shows that each $\alpha$-chain has at most four roots. Thus the possible positions of $\beta$ in its $\alpha$-chain are


In the first case we have $\langle\alpha+\beta, \alpha+\beta\rangle=\langle\beta, \beta\rangle$ since $s_{\alpha}(\beta)=\alpha+\beta$. In the remaining cases the first and last roots in the $\alpha$-chain are long roots and the remainder are short roots. The relative lengths are given in the proof of Proposition 6.1. We have

$$
\frac{\langle\alpha+\beta, \alpha+\beta\rangle}{\langle\beta, \beta\rangle}=1, \frac{1}{2}, 2, \frac{1}{3}, 1,3
$$

in the above six cases respectively. Thus in each case we have

$$
\frac{\langle\alpha+\beta, \alpha+\beta\rangle}{\langle\beta, \beta\rangle}=\frac{p+1}{q}
$$

and so $N_{\alpha, \beta} N_{-\alpha,-\beta}=-(p+1)^{2}$
(iv) Now suppose that $\alpha, \beta, \gamma, \delta \in \Phi$ satisfy $\alpha+\beta+\gamma+\delta=0$ with no pair equal and opposite. Consider the Jacobi identity

$$
\left[\left[e_{\alpha} e_{\beta}\right] e_{\gamma}\right]+\left[\left[e_{\beta} e_{\gamma}\right] e_{\alpha}\right]+\left[\left[e_{\gamma} e_{\alpha}\right] e_{\beta}\right]=0
$$

This gives

$$
N_{\alpha, \beta} N_{\alpha+\beta, \gamma}+N_{\beta, \gamma} N_{\beta+\gamma, \alpha}+N_{\gamma, \alpha} N_{\gamma+\alpha, \beta}=0
$$

Using relations (ii) this gives

$$
\frac{N_{\alpha, \beta} N_{\gamma, \delta}}{\langle\alpha+\beta, \alpha+\beta\rangle}+\frac{N_{\beta, \gamma} N_{\alpha, \delta}}{\langle\beta+\gamma, \beta+\gamma\rangle}+\frac{N_{\gamma, \alpha} N_{\beta, \delta}}{\langle\gamma+\alpha, \gamma+\alpha\rangle}=0 .
$$

(As usual we interpret $N_{\theta, \phi}$ as 0 if $\theta+\phi$ is not a root.)

Proposition 7.1 (iii) has a very useful corollary.

Corollary 7.2 If $\alpha, \beta, \alpha+\beta \in \Phi$ then $N_{\alpha, \beta} \neq 0$. Thus $\left[L_{\alpha} L_{\beta}\right]=L_{\alpha+\beta}$.

### 7.2 The uniqueness theorem

We shall now use the above relations between the structure constants to show that the Lie algebra $L$ is uniquely determined up to isomorphism. A Dynkin diagram on the standard list 6.11 is given, and this determines uniquely a Cartan matrix $A=\left(A_{i j}\right)$ on the standard list 6.12. Now the Cartan matrix determines the set $\Phi$ of roots as linear combinations of the fundamental roots $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$. For each root $\alpha \in \Phi$ has form $\alpha=w\left(\alpha_{i}\right)$ for some $\alpha_{i} \in \Pi$ and some $w \in W$, by Proposition 5.12. Moreover each element $w \in W$ is a product of elements $s_{1}, \ldots, s_{l}$ by Theorem 5.13. The actions of $s_{1}, \ldots, s_{l}$ on the fundamental roots $\alpha_{1}, \ldots, \alpha_{l}$ are given in terms of the Cartan matrix by

$$
s_{i}\left(\alpha_{j}\right)=\alpha_{j}-A_{i j} \alpha_{i}
$$

Thus by applying the fundamental reflections successively to the fundamental roots we obtain all roots as linear combinations of the fundamental roots.

We next observe that all scalar products $\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle$ for $\alpha, \beta \in \Phi$ are determined by the Cartan matrix. By Proposition $4.222 \frac{\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}$ is determined by the root system, hence by the Cartan matrix as shown above. Then $\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle$ is determined by the formula

$$
\frac{1}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}=\sum_{\beta \in \Phi}\left(\frac{\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}\right)^{2}
$$

of Proposition 4.24. Thus $\left\langle h_{\alpha}^{\prime}, h_{\beta}^{\prime}\right\rangle$ is also determined by the Cartan matrix. Thus we see that if the structure constants $N_{\alpha, \beta}$ are known the multiplication of basis elements

$$
\begin{aligned}
& {\left[h_{i} h_{j}\right]=0} \\
& {\left[h_{i} e_{\alpha}\right]=\alpha\left(h_{i}\right) e_{\alpha}} \\
& {\left[e_{\alpha} e_{-\alpha}\right]=h_{\alpha}} \\
& {\left[e_{\alpha} e_{\beta}\right]=N_{\alpha, \beta} e_{\alpha+\beta} \quad \text { if } \alpha, \beta, \alpha+\beta \in \Phi} \\
& {\left[e_{\alpha} e_{\beta}\right]=0 \quad \text { if } \alpha+\beta \neq 0 \text { and } \alpha+\beta \notin \Phi}
\end{aligned}
$$

will be completely determined.
We shall show that for certain pairs $(\alpha, \beta)$ of roots the structure constants $N_{\alpha, \beta}$ can be chosen arbitrarily, and that the remaining structure constants are uniquely determined in terms of these by the relations of Proposition 7.1.

We choose a total ordering on the vector space $V=H_{\mathbb{R}}^{*}$ as in Section 5.1 giving rise to the positive system $\Phi^{+}$and fundamental system $\Pi$ of roots. An ordered pair $(\alpha, \beta)$ of roots will be called special if $\alpha+\beta \in \Phi$ and $0<\alpha<\beta$. The pair $(\alpha, \beta)$ will be called extraspecial if $(\alpha, \beta)$ is special and if, in addition, for all special pairs $(\gamma, \delta)$ such that $\alpha+\beta=\gamma+\delta$ we have $\alpha \leq \gamma$.

Lemma 7.3 The structure constants $N_{\alpha, \beta}$ for extraspecial pairs $(\alpha, \beta)$ can be chosen as arbitrary non-zero elements of $\mathbb{C}$, by appropriate choice of the elements $e_{\alpha}$.

Proof. We choose the $e_{\alpha}$ for $\alpha \in \Phi^{+}$in the order given by $<$. Suppose $(\alpha, \beta)$ is an extraspecial pair. Then we have

$$
\left[e_{\alpha} e_{\beta}\right]=N_{\alpha, \beta} e_{\alpha+\beta}
$$

and $e_{\alpha}, e_{\beta}$ have already been chosen. Moreover there is only one extraspecial pair with given sum $\alpha+\beta$. Thus $e_{\alpha+\beta}$ can be chosen to give any non-zero value of $N_{\alpha, \beta}$.

Proposition 7.4 All the structure constants $N_{\alpha, \beta}$ are determined by the structure constants for extraspecial pairs.

Proof. We consider the set of all pairs of roots $(\alpha, \beta)$ such that $\alpha+\beta$ is a root. Let $(\alpha, \beta)$ be such a pair and let $\gamma=-\alpha-\beta$. Then the following 12 pairs of roots are of the given type.

$$
\begin{aligned}
& (\alpha, \beta)(\beta, \gamma)(\gamma, \alpha)(\beta, \alpha)(\gamma, \beta)(\alpha, \gamma) \\
& (-\alpha,-\beta)(-\beta,-\gamma)(-\gamma,-\alpha)(-\beta,-\alpha)(-\gamma,-\beta)(-\alpha,-\gamma)
\end{aligned}
$$

Since $\alpha+\beta+\gamma=0$ either two or one of $\alpha, \beta, \gamma$ are positive. Thus either two of $\alpha, \beta, \gamma$ are positive or two of $-\alpha,-\beta,-\gamma$ are positive. By choosing two positive roots from $\alpha, \beta, \gamma$ or from $-\alpha,-\beta,-\gamma$ and by writing them in the appropriate order we obtain a special pair. Thus just one of the above 12 pairs of roots is a special pair.

Now the relations in Proposition 7.1 (i), (ii), (iii) enable us to express $N_{\beta, \alpha}, N_{\beta, \gamma}, N_{\gamma, \alpha}$ and $N_{-\alpha,-\beta}$ in terms of $N_{\alpha, \beta}$. Thus these relations enable us to express $N_{\theta, \phi}$ for all the 12 pairs $(\theta, \phi)$ above in terms of $N_{\theta, \phi}$ for the special pair $(\theta, \phi)$.

The next stage is to show that the $N_{\alpha, \beta}$ for all special pairs $(\alpha, \beta)$ are determined in terms of the $N_{\alpha, \beta}$ for extraspecial pairs. Suppose $(\alpha, \beta)$ is special but not extraspecial. Then there exists an extraspecial pair $(\gamma, \delta)$ such that $\alpha+\beta=\gamma+\delta$. Thus $\alpha+\beta+(-\gamma)+(-\delta)=0$ and no pair of $\alpha, \beta,-\gamma,-\delta$ are equal and opposite. By Proposition 7.1 (iv) we have

$$
\frac{N_{\alpha, \beta} N_{-\gamma,-\delta}}{\langle\alpha+\beta, \alpha+\beta\rangle}+\frac{N_{\beta,-\gamma} N_{\alpha,-\delta}}{\langle\beta-\gamma, \beta-\gamma\rangle}+\frac{N_{-\gamma, \alpha} N_{\beta,-\delta}}{\langle-\gamma+\alpha,-\gamma+\alpha\rangle}=0 .
$$

Now the roots $\alpha, \beta, \gamma, \delta$ are ordered by

$$
0<\gamma<\alpha<\beta<\delta
$$

Thus we may use relations (i), (ii), (iii) of Proposition 7.1 to express $N_{-\gamma,-\delta}$ in terms of $N_{\gamma, \delta} ; N_{\beta,-\gamma}$ in terms of $N_{\gamma, \beta-\gamma} ; N_{\alpha,-\delta}$ in terms of $N_{\alpha, \delta-\alpha} ; N_{-\gamma, \alpha}$ in terms of $N_{\gamma, \alpha-\gamma}$; and $N_{\beta,-\delta}$ in terms of $N_{\beta, \delta-\beta}$. Thus $N_{\alpha, \beta}$ is expressed in terms of

$$
N_{\gamma, \delta}, N_{\gamma, \beta-\gamma}, N_{\alpha, \delta-\alpha}, N_{\gamma, \alpha-\gamma}, N_{\beta, \delta-\beta}
$$

Now $(\gamma, \delta)$ is an extraspecial pair and $(\gamma, \beta-\gamma),(\alpha, \delta-\alpha),(\gamma, \alpha-\gamma)$ and $(\beta, \delta-\beta)$ are all pairs of positive roots whose sums are roots less than $\alpha+\beta=\gamma+\delta$ in the given ordering. We may therefore argue by induction on $\alpha+\beta$, using the given order, that $N_{\alpha, \beta}$ can be expressed in terms of $N_{\theta, \phi}$ for extraspecial pairs $(\theta, \phi)$.

We can now state our uniqueness theorem.

Theorem 7.5 Any two simple Lie algebras with the same Cartan matrix are isomorphic.

Proof. We choose the basis elements $\left\{h_{i}, e_{\alpha}\right\}$ of such a Lie algebra $L$ such that $N_{\alpha, \beta}=1$ for all extraspecial pairs of roots $(\alpha, \beta)$. We may do this by Lemma 7.3. The remaining structure constants $N_{\alpha, \beta}$ are all then uniquely determined by Proposition 7.4. Thus the formulae expressing a Lie product of
basis elements as a linear combination of basis elements are completely determined by the Cartan matrix. Thus the Lie algebra $L$ is uniquely determined up to isomorphism.

### 7.3 Some generators and relations in a simple Lie algebra

We now turn to the question of the existence of a simple Lie algebra with Cartan matrix on the standard list 6.12. A proof of the existence theorem has been given by J. Tits (IHES Publ. Math. 31 (1966)) along the lines of the arguments used so far. The details are technically quite complicated, however, and so we prefer to give a different proof of the existence theorem.

Let $L$ be a simple Lie algebra with Cartan matrix $A$. Let $H$ be a Cartan subalgebra of $L$ and

$$
L=H \oplus \sum_{\alpha \in \Phi} L_{\alpha}
$$

be the Cartan decomposition. As before we consider the elements $h_{i} \in H$ given by

$$
h_{i}=\frac{2 h_{\alpha_{i}}^{\prime}}{\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle}
$$

where $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$ is a fundamental system in $\Phi$. As in Section 7.1 we can choose elements $e_{i} \in L_{\alpha_{i}}, f_{i} \in L_{-\alpha_{i}}$ such that $\left[e_{i} f_{i}\right]=h_{i}$.

We shall show that the elements $e_{1}, \ldots, e_{l}, h_{1}, \ldots, h_{l}, f_{1}, \ldots, f_{l}$ generate $L$. (Of course this is equivalent to saying that $e_{1}, \ldots, e_{l}, f_{1}, \ldots, f_{l}$ generate $L$, but it will be useful to include $h_{1}, \ldots, h_{l}$ in the generating set.)

Lemma 7.6 If $\alpha \in \Phi^{+}$and $\alpha \notin \Pi$ there exists $\alpha_{i} \in \Pi$ such that $\alpha-\alpha_{i} \in \Phi^{+}$. Thus every positive non-fundamental root is the sum of a fundamental root with a positive root.

Proof. Suppose if possible that the result is false. Then $\alpha-\alpha_{i}$ is not a root and is non-zero for each $i$. (We can use Corollary 5.6 to see that $\alpha-\alpha_{i}$ cannot be a negative root.) Consider the $\alpha_{i}$-chain of roots through $\alpha$. This has form

$$
\alpha, \alpha_{i}+\alpha, \ldots, q \alpha_{i}+\alpha
$$

By Proposition 4.22 we have

$$
2 \frac{\left\langle\alpha_{i}, \alpha\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}=-q
$$

This implies that $\left\langle\alpha_{i}, \alpha\right\rangle \leq 0$. Now $\alpha \in \Phi^{+}$has form $\alpha=\sum_{i} n_{i} \alpha_{i}$ with all $n_{i} \geq 0$. Thus

$$
\langle\alpha, \alpha\rangle=\sum_{i} n_{i}\left\langle\alpha_{i}, \alpha\right\rangle \leq 0
$$

This gives a contradiction, since we know $\langle\alpha, \alpha\rangle>0$.
Proposition 7.7 The elements $e_{1}, \ldots, e_{l}, h_{1}, \ldots, h_{l}, f_{1}, \ldots, f_{l}$ generate $L$.

Proof. Since $h_{1}, \ldots, h_{l}$ span $H$ it will be sufficient to show that each $L_{\alpha}$ for $\alpha \in \Phi^{+}$lies in the subalgebra generated by $e_{1}, \ldots, e_{l}$ and each $L_{\alpha}$ for $\alpha \in \Phi^{-}$ lies in the subalgebra generated by $f_{1}, \ldots, f_{l}$.

Let $\alpha \in \Phi^{+}$. If $\alpha=\alpha_{i}$ for some $i$ we have $L_{\alpha}=\mathbb{C} e_{i}$. If $\alpha \notin \Pi$ we can write $\alpha=\alpha_{i}+\beta$ for some $\alpha_{i} \in \Pi$ and some $\beta \in \Phi^{+}$by Lemma 7.6. We then have $\left[L_{\alpha_{i}} L_{\beta}\right]=L_{\alpha}$ by Corollary 7.2. Thus we may choose $e_{\alpha}=\left[e_{i}, e_{\beta}\right]$ for some $e_{\beta} \neq 0$ in $L_{\beta}$. By repeating this process we obtain

$$
e_{\alpha}=\left[\left[e_{i_{1}} e_{i_{2}}\right] \ldots e_{i_{k}}\right]
$$

for some sequence $i_{1}, \ldots, i_{k}$. Thus each $L_{\alpha}$ for $\alpha \in \Phi^{+}$lies in the subalgebra generated by $e_{1}, \ldots, e_{l}$. Similarly each $L_{\alpha}$ for $\alpha \in \Phi^{-}$lies in the subalgebra generated by $f_{1}, \ldots, f_{l}$.

Proposition 7.8 The generators $e_{1}, \ldots, e_{l}, h_{1}, \ldots, h_{l}, f_{1}, \ldots, f_{l}$ of $L$ satisfy the following relations.
(a) $\left[h_{i} h_{j}\right]=0$
(b) $\left[h_{i} e_{j}\right]=A_{i j} e_{j}$
(c) $\left[h_{i} f_{j}\right]=-A_{i j} f_{j}$
(d) $\left[e_{i} f_{i}\right]=h_{i}$
(e) $\left[e_{i} f_{j}\right]=0 \quad$ if $i \neq j$
(f) $\left[e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right]=0 \quad$ if $i \neq j$
$\leftarrow 1-A_{i j} \rightarrow$
(g) $\left[f_{i}\left[f_{i} \ldots\left[f_{i} f_{j}\right]\right]\right]=0 \quad$ if $i \neq j$. $\leftarrow 1-A_{i j} \rightarrow$
Note that in relations $(\mathrm{f}),(\mathrm{g})$ there are $1-A_{i j}$ occurrences of $e_{i}, f_{i}$ respectively. Since $A_{i j} \leq 0$ for $i \neq j$ this number $1-A_{i j}$ is a positive integer.

Proof. Relation (a) follows from $[H H]=0$. For relation (b), we have

$$
\begin{aligned}
{\left[h_{i} e_{j}\right] } & =2 \frac{\left[h_{\alpha_{i}}^{\prime} e_{j}\right]}{\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle}=2 \frac{\alpha_{j}\left(h_{\alpha_{i}}^{\prime}\right)}{\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle} e_{j} \\
& =2 \frac{\left\langle h_{\alpha_{j}}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle}{\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle} e_{j}=2 \frac{\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} e_{j}=A_{i j} e_{j}
\end{aligned}
$$

Relation (c) is obtained similarly. Relation (d) holds by definition of $f_{i}$. Relation (e) holds because $\left[e_{i} f_{j}\right] \in L_{\alpha_{i}-\alpha_{j}}$ and $\alpha_{i}-\alpha_{j}$ is not a root when $i \neq j$, as follows from Corollary 5.6. In order to prove relation (f) we consider the $\alpha_{i}$-chain of roots through $\alpha_{j}$. Since $-\alpha_{i}+\alpha_{j}$ is not a root this chain has form

$$
\alpha_{j}, \alpha_{i}+\alpha_{j}, \ldots, q \alpha_{i}+\alpha_{j}
$$

By Proposition 4.22 we have $A_{i j}=-q$. Thus $\left(1-A_{i j}\right) \alpha_{i}+\alpha_{j}$ is not a root. Since the element $\left[e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right]$ lies in $L_{\left(1-A_{i j}\right) \alpha_{i}+\alpha_{j}}$ this element must be 0 . Relation (g) is obtained similarly.

### 7.4 The Lie algebras $L(A)$ and $\tilde{L}(A)$

Let $A$ be a Cartan matrix on the standard list 6.12 . Motivated by Propositions 7.7 and 7.8 we shall construct a Lie algebra $L(A)$ which will be shown to be a finite dimensional simple Lie algebra with Cartan matrix $A$.

Suppose $A$ is an $l \times l$ matrix. Let $\mathfrak{F}$ be the free associative algebra over $\mathbb{C}$ on the $3 l$ generators $e_{1}, \ldots, e_{l}, h_{1}, \ldots, h_{l}, f_{1}, \ldots, f_{l}$. The set of all monomials in these generators form a basis for $\mathfrak{F}$. Let [ $\mathfrak{F}$ ] be the Lie algebra obtained from $\mathfrak{F}$ by redefining the multiplication in the usual way and let $\mathfrak{Z}$ be the subalgebra of $[\mathfrak{F}]$ generated by the elements $e_{1}, \ldots, e_{l}, h_{1}, \ldots, h_{l}, f_{1}, \ldots, f_{l}$. Let $J$ be the ideal of $\mathfrak{R}$ generated by the elements

$$
\begin{aligned}
& {\left[h_{i} h_{j}\right]} \\
& {\left[h_{i} e_{j}\right]-A_{i j} e_{j}} \\
& {\left[h_{i} f_{j}\right]+A_{i j} f_{j}} \\
& {\left[e_{i} f_{i}\right]-h_{i}} \\
& {\left[e_{i} f_{j}\right] \quad \text { for } i \neq j} \\
& {\left[e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right]} \\
& {\left[f_{i}\left[f_{i} \ldots\left[f_{i} f_{j}\right]\right]\right] \quad \text { for } i \neq j} \\
& \text { for } i \neq j
\end{aligned}
$$

where the number of occurrences of $e_{i}, f_{i}$ respectively in the last two elements is $1-A_{i j}$.

We define $L(A)=\mathbb{R} / J$. We shall eventually be able to show that $L(A)$ is the Lie algebra we require to prove the existence theorem. This description of $L(A)$ by generators and relations is due to J. P. Serre.

In order to investigate the Lie algebra $L(A)$ it is convenient to define a second, larger, Lie algebra $\tilde{L}(A)$. Let $\tilde{J}$ be the ideal of $\mathcal{R}$ generated by the elements

$$
\begin{aligned}
& {\left[h_{i} h_{j}\right]} \\
& {\left[h_{i} e_{j}\right]-A_{i j} e_{j}} \\
& {\left[h_{i} f_{j}\right]+A_{i j} f_{j}} \\
& {\left[e_{i} f_{i}\right]-h_{i}} \\
& {\left[e_{i} f_{j}\right] \quad \text { for } i \neq j}
\end{aligned}
$$

Let $\tilde{L}(A)=\mathbb{Z} / \tilde{J}$. Since $\tilde{J} \subset J$ we have surjective Lie algebra homomorphisms

$$
\mathfrak{Z} \rightarrow \tilde{L}(A) \rightarrow L(A)
$$

We shall investigate the properties of the Lie algebra $\tilde{L}(A)$. This is generated by the images of the generators of $\mathfrak{R}$ under the above homomorphism. These images will continue to be written $e_{1}, \ldots, e_{l}, h_{1}, \ldots, h_{l}, f_{1}, \ldots, f_{l}$. These elements satisfy the relations

$$
\begin{aligned}
& {\left[h_{i} h_{j}\right]=0} \\
& {\left[h_{i} e_{j}\right]=A_{i j} e_{j}} \\
& {\left[h_{i} f_{j}\right]=-A_{i j} f_{j}} \\
& {\left[e_{i} f_{i}\right]=h_{i}} \\
& {\left[e_{i} f_{j}\right]=0 \quad \text { for } i \neq j}
\end{aligned}
$$

Proposition 7.9 Let $\mathfrak{F}^{-}$be the free associative algebra over $\mathbb{C}$ with generators $f_{1}, \ldots, f_{l}$. Then $\mathfrak{F}^{-}$may be made into an $\tilde{L}(A)$-module giving a representation $\rho: \tilde{L}(A) \rightarrow\left[\right.$ End $\left.\mathfrak{F}^{-}\right]$defined by:

$$
\begin{aligned}
& \rho\left(f_{i}\right) f_{i_{1}} \ldots f_{i_{r}}=f_{i} f_{i_{1}} \ldots f_{i_{r}} \\
& \rho\left(h_{i}\right) f_{i_{1}} \ldots f_{i_{r}}=-\left(\sum_{k=1}^{r} A_{i i_{k}}\right) f_{i_{1}} \ldots f_{i_{r}} \\
& \rho\left(e_{i}\right) f_{i_{1}} \ldots f_{i_{r}}=-\sum_{k=1}^{r} \delta_{i i_{k}}\left(\sum_{h=k+1}^{r} A_{i i_{h}}\right) f_{i_{1}} \ldots \hat{f}_{i_{k}} \ldots f_{i_{r}}
\end{aligned}
$$

where as usual the symbol $\hat{f}_{i_{k}}$ means that $f_{i_{k}}$ is omitted from the product.

Proof. Since the monomials $f_{i_{1}} \ldots f_{i_{r}}$ form a basis for $\mathfrak{F}^{-}$the endomorphisms $\rho\left(f_{i}\right), \rho\left(h_{i}\right), \rho\left(e_{i}\right)$ are uniquely determined by the above formulae. Thus there is a unique homomorphism $\mathfrak{F} \rightarrow$ End $\mathfrak{F}^{-}$mapping $e_{i}, h_{i}, f_{i}$ to $\rho\left(e_{i}\right), \rho\left(h_{i}\right), \rho\left(f_{i}\right)$ respectively. This induces a Lie algebra homomorphism $[\mathfrak{F}] \rightarrow\left[\right.$ End $\left.\mathfrak{F}^{-}\right]$and so, by restriction, a Lie algebra homomorphism $\mathbb{R} \rightarrow$ [End $\mathscr{F}^{-}$]. In order to obtain a homomorphism $\tilde{L}(A) \rightarrow\left[\right.$ End $\left.\mathfrak{F}^{-}\right]$we must verify the following relations.
(a) $\left[\rho\left(h_{i}\right) \rho\left(h_{j}\right)\right]=0$
(b) $\left[\rho\left(h_{i}\right) \rho\left(e_{j}\right)\right]=A_{i j} \rho\left(e_{j}\right)$
(c) $\left[\rho\left(h_{i}\right) \rho\left(f_{j}\right)\right]=-A_{i j} \rho\left(f_{j}\right)$
(d) $\left[\rho\left(e_{i}\right) \rho\left(f_{i}\right)\right]=\rho\left(h_{i}\right)$
(e) $\left[\rho\left(e_{i}\right) \rho\left(f_{j}\right)\right]=0 \quad$ for $i \neq j$

Relation (a) is trivial since $\rho\left(h_{i}\right)$ multiplies each basis element of $\mathfrak{F}^{-}$by a scalar.

To prove relation (b) we have

$$
\begin{aligned}
& \rho\left(h_{i}\right) \rho\left(e_{j}\right) f_{i_{1}} \ldots f_{i_{r}}=-\sum_{k=1}^{r} \delta_{j i_{k}}\left(\sum_{h=k+1}^{r} A_{j i_{h}}\right)\left(-\sum_{g \neq k} A_{i i_{g}}\right) f_{i_{1}} \ldots \hat{f}_{i_{k}} \ldots f_{i_{r}} \\
& \rho\left(e_{j}\right) \rho\left(h_{i}\right) f_{i_{1}} \ldots f_{i_{r}}=-\sum_{k=1}^{r} \delta_{j i_{k}}\left(\sum_{h=k+1}^{r} A_{j i_{h}}\right)\left(-\sum_{g=1}^{r} A_{i i_{g}}\right) f_{i_{l}} \ldots \hat{f}_{i_{k}} \ldots f_{i_{r}}
\end{aligned}
$$

Thus

$$
\begin{aligned}
& \left(\rho\left(h_{i}\right) \rho\left(e_{j}\right)-\rho\left(e_{j}\right) \rho\left(h_{i}\right)\right) f_{i_{1}} \ldots f_{i_{r}} \\
& \quad=A_{i j}\left(-\sum_{k=1}^{r} \delta_{j i_{k}}\left(\sum_{h=k+1}^{r} A_{j i_{h}}\right)\right) f_{i_{1}} \ldots \hat{f}_{i_{k}} \ldots f_{i_{r}} \\
& \quad=A_{i j} \rho\left(e_{j}\right) f_{i_{1}} \ldots f_{i_{r}}
\end{aligned}
$$

To prove relation (c) we have

$$
\begin{aligned}
& \rho\left(h_{i}\right) \rho\left(f_{j}\right) f_{i_{1}} \ldots f_{i_{r}}=-\left(A_{i j}+\sum_{k=1}^{r} A_{i i_{k}}\right) f_{j} f_{i_{1}} \ldots f_{i_{r}} \\
& \rho\left(f_{j}\right) \rho\left(h_{i}\right) f_{i_{1}} \ldots f_{i_{r}}=-\left(\sum_{k=1}^{r} A_{i i_{k}}\right) f_{j} f_{i_{1}} \ldots f_{i_{r}} .
\end{aligned}
$$

Thus

$$
\left(\rho\left(h_{i}\right) \rho\left(f_{j}\right)-\rho\left(f_{j}\right) \rho\left(h_{i}\right)\right) f_{i_{1}} \ldots f_{i_{r}}=-A_{i j} f_{j} f_{i_{1}} \ldots f_{i_{r}}=-A_{i j} \rho\left(f_{j}\right) f_{i_{1}} \ldots f_{i_{r}}
$$

We next consider relation (d). We have

$$
\begin{aligned}
\rho\left(e_{i}\right) \rho\left(f_{i}\right) f_{i_{1}} \ldots f_{i_{r}}= & -\left(\sum_{h=1}^{r} A_{i i_{h}}\right) f_{i_{1}} \ldots f_{i_{r}} \\
& -\sum_{k=1}^{r} \delta_{i i_{k}}\left(\sum_{h=k+1}^{r} A_{i i_{h}}\right) f_{i} f_{i_{1}} \ldots \hat{f}_{i_{k}} \ldots f_{i_{r}} \\
\rho\left(f_{i}\right) \rho\left(e_{i}\right) f_{i_{1}} \ldots f_{i_{r}}= & -\sum_{k=1}^{r} \delta_{i i_{k}}\left(\sum_{h=k+1}^{r} A_{i i_{h}}\right) f_{i} f_{i_{1}} \ldots \hat{f}_{i_{k}} \ldots f_{i_{r}} .
\end{aligned}
$$

Thus

$$
\left(\rho\left(e_{i}\right) \rho\left(f_{i}\right)-\rho\left(f_{i}\right) \rho\left(e_{i}\right)\right) f_{i_{1}} \ldots f_{i_{r}}=-\left(\sum_{h=1}^{r} A_{i i_{h}}\right) f_{i_{1}} \ldots f_{i_{r}}=\rho\left(h_{i}\right) f_{i_{1}} \ldots f_{i_{r}}
$$

Finally we consider relation (e). Suppose $i \neq j$. Then

$$
\begin{aligned}
\rho\left(e_{i}\right) \rho\left(f_{j}\right) f_{i_{1}} \ldots f_{i_{r}} & =-\sum_{k=1}^{r} \delta_{i i_{k}}\left(\sum_{h=k+1}^{r} A_{i i_{h}}\right) f_{j} f_{i_{1}} \ldots \hat{f}_{i_{k}} \ldots f_{i_{r}} \\
& =\rho\left(f_{j}\right) \rho\left(e_{i}\right) f_{i_{1}} \ldots f_{i_{r}}
\end{aligned}
$$

Thus all the relations are preserved and we have a homomorphism $\tilde{L}(A) \rightarrow$ [End $\mathfrak{F}^{-}$].

We can deduce useful information about $\tilde{L}(A)$ from the existence of this homomorphism.

Proposition 7.10 The elements $h_{1}, \ldots, h_{l}$ of $\tilde{L}(A)$ are linearly independent.

Proof. We show that the elements $\rho\left(h_{1}\right), \ldots, \rho\left(h_{l}\right)$ of End $\mathfrak{F}^{-}$are linearly independent. We have

$$
\rho\left(h_{i}\right) f_{j}=-A_{i j} f_{j} .
$$

Thus if $\sum \lambda_{i} \rho\left(h_{i}\right)=0$ we would have $\sum_{i} \lambda_{i} A_{i j}=0$ for all $j=1, \ldots, l$. Since the Cartan matrix $A=\left(A_{i j}\right)$ is non-singular this implies that $\lambda_{i}=0$ for each $i$. Hence $\rho\left(h_{1}\right), \ldots, \rho\left(h_{l}\right)$ are linearly independent, and so $h_{1}, \ldots, h_{l}$ must be linearly independent also.

Let $\tilde{H}$ be the subspace of $\tilde{L}(A)$ spanned by $h_{1}, \ldots, h_{l}$. Then we have $\operatorname{dim} \tilde{H}=l$. Moreover $[\tilde{H} \tilde{H}]=O$, thus $\tilde{H}$ is an abelian subalgebra of $\tilde{L}(A)$. We consider the weight spaces of $\tilde{L}(A)$ with respect to $\tilde{H}$. We are no longer dealing with a finite dimensional $\tilde{H}$-module as in Theorem 2.9, but analogous
ideas apply in our situation. Elements of $\operatorname{Hom}(\tilde{H}, \mathbb{C})$ will be called weights. For each weight $\mu: \tilde{H} \rightarrow \mathbb{C}$ we define the corresponding weight space $\tilde{L}(A)_{\mu}$ by

$$
\tilde{L}(A)_{\mu}=\{x \in \tilde{L}(A) ;[h x]=\mu(h) x \quad \text { for all } h \in \tilde{H}\} .
$$

Proposition 7.11 $\tilde{L}(A)=\bigoplus_{\mu} \tilde{L}(A)_{\mu}$. Thus $\tilde{L}(A)$ is the direct sum of its weight spaces.

Proof. We first show that $\tilde{L}(A)=\sum_{\mu} \tilde{L}(A)_{\mu}$. A vector which lies in a weight space will be called a weight vector. We observe that, if $x, y \in \tilde{L}(A)$ are weight vectors of weights $\lambda, \mu$ respectively, then $[x y]$ is a weight vector of weight $\lambda+\mu$. For we have

$$
\begin{aligned}
{[h[x y]] } & =[[h x] y]+[x[h y]]=\lambda(h)[x y]+\mu(h)[x y] \\
& =(\lambda+\mu)(h)[x y] \quad \text { for } h \in \tilde{H} .
\end{aligned}
$$

Now $\tilde{L}(A)$ is generated by elements $e_{i}, h_{i}, f_{i}$. Let $\alpha_{i} \in \operatorname{Hom}(\tilde{H}, \mathbb{C})$ be defined by

$$
\alpha_{i}\left(h_{j}\right)=A_{j i} .
$$

Then $e_{i}$ is a weight vector of weight $\alpha_{i}, f_{i}$ is a weight vector of weight $-\alpha_{i}$ and $h_{i}$ is a weight vector of weight 0 . Thus all Lie products of generators $e_{i}$, $h_{i}, f_{i}$ are weight vectors. Since every element of $\tilde{L}(A)$ is a linear combination of such products we deduce that

$$
\tilde{L}(A)=\sum_{\mu} \tilde{L}(A)_{\mu} .
$$

We next show that this sum is direct. If this is not so we can find a non-zero vector $x \in \tilde{L}(A)_{\mu}$ such that $x=\sum_{\nu} x_{\nu}$ where $x_{\nu} \in \tilde{L}(A)_{\nu}$ and $\nu$ runs over a finite set of weights all distinct from $\mu$. Since $x \in \tilde{L}(A)_{\mu}$ we have

$$
(\operatorname{ad} h-\mu(h) 1) x=0 .
$$

Since $x=\sum_{\nu} x_{\nu}$ with $x_{\nu} \in \tilde{L}(A)_{\nu}$ we have

$$
\prod_{\nu}(\operatorname{ad} h-\nu(h) 1) x=0 .
$$

Now we can find an element $h \in \tilde{H}$ such that $\mu(h) \neq \nu(h)$ for all such $\nu$. For the elements satisfying $\mu(h)=\nu(h)$ for some fixed $\nu$ lie in a proper subspace of $\tilde{H}$, and the finite dimensional vector space $\tilde{H}$ over $\mathbb{C}$ cannot be expressed
as the union of a finite number of proper subspaces. Thus we choose $h \in \tilde{H}$ such that $\mu(h) \neq \nu(h)$ for all such $\nu$. Then the polynomials

$$
t-\mu(h), \quad \prod_{\nu}(t-\nu(h))
$$

in $\mathbb{C}[t]$ are coprime. Thus there exist polynomials $a(t), b(t) \in \mathbb{C}[t]$ with

$$
a(t)(t-\mu(h))+b(t) \prod_{\nu}(t-\nu(h))=1
$$

If follows that

$$
a(\operatorname{ad} h)(\operatorname{ad} h-\mu(h) 1) x+b(\operatorname{ad} h) \prod_{\nu}(\operatorname{ad} h-\nu(h) 1) x=x .
$$

We deduce that $x=0$, a contradiction. Thus the $\operatorname{sum} \sum_{\mu} \tilde{L}(A)_{\mu}$ is direct.
We next obtain information about the kind of weights $\mu$ which can occur, that is for which $\tilde{L}(A)_{\mu} \neq 0$. The weights $\alpha_{1}, \ldots, \alpha_{l} \in \operatorname{Hom}(\tilde{H}, \mathbb{C})$ are linearly independent since the Cartan matrix $A$ is non-singular. Thus any weight has form $n_{1} \alpha_{1}+\cdots+n_{l} \alpha_{l}$ for $n_{i} \in \mathbb{C}$. We shall show that all weights $\mu$ which occur in $\tilde{L}(A)$ have this form with $n_{i} \in \mathbb{Z}$ and with either $n_{i} \geq 0$ for all $i$ or $n_{i} \leq 0$ for all $i$.

Let

$$
\begin{aligned}
& Q=\left\{n_{1} \alpha_{1}+\cdots+n_{l} \alpha_{l} ; n_{i} \in \mathbb{Z}\right\} \\
& Q^{+}=\left\{n_{1} \alpha_{1}+\cdots+n_{l} \alpha_{l} \neq 0 ; n_{i} \geq 0 \quad \text { for all } i\right\} \\
& Q^{-}=\left\{n_{1} \alpha_{1}+\cdots+n_{l} \alpha_{l} \neq 0 ; n_{i} \leq 0 \quad \text { for all } i\right\}
\end{aligned}
$$

Let

$$
\begin{aligned}
& \tilde{L}(A)^{+}=\sum_{\mu \in Q^{+}} \tilde{L}(A)_{\mu} \\
& \tilde{L}(A)^{-}=\sum_{\mu \in Q^{-}} \tilde{L}(A)_{\mu}
\end{aligned}
$$

It follows from Proposition 7.11 that the sum $\tilde{L}(A)^{-}+\tilde{H}+\tilde{L}(A)^{+}$is direct. We shall show that in fact

$$
\tilde{L}(A)=\tilde{L}(A)^{-} \oplus \tilde{H} \oplus \tilde{L}(A)^{+}
$$

Let $\tilde{N}$ be the subalgebra of $\tilde{L}(A)$ generated by $e_{1}, \ldots, e_{l}$ and $\tilde{N}^{-}$the subalgebra generated by $f_{1}, \ldots, f_{l}$. Since $e_{i}$ has weight $\alpha_{i}$ and $f_{i}$ has weight $-\alpha_{i}$ we have $\tilde{N} \subset \tilde{L}(A)^{+}$and $\tilde{N}^{-} \subset \tilde{L}(A)^{-}$. Thus the sum $\tilde{N}^{-}+\tilde{H}+\tilde{N}$ is direct.

Proposition 7.12 (i) $\tilde{L}(A)=\tilde{N}^{-} \oplus \tilde{H} \oplus \tilde{N}$
(ii) $\tilde{N}=\tilde{L}(A)^{+}, \quad \tilde{N}^{-}=\tilde{L}(A)^{-}, \quad \tilde{H}=\tilde{L}(A)_{0}$
(iii) Every non-zero weight of $\tilde{L}(A)$ lies in $Q^{+}$or in $Q^{-}$.

Proof. The relations $\left[h_{i} e_{j}\right]=A_{i j} e_{j}$ show that $\left[h_{i}, \tilde{N}\right] \subset \tilde{N}$ since the $e_{j}$ generate $\tilde{N}$. Thus we have $[\tilde{H}, \tilde{N}] \subset \tilde{N}$. It follows that $\tilde{H}+\tilde{N}$ is a subalgebra of $\tilde{L}(A)$, since

$$
[\tilde{H}+\tilde{N}, \tilde{H}+\tilde{N}] \subset[\tilde{H} \tilde{H}]+[\tilde{H} \tilde{N}]+[\tilde{N} \tilde{N}] \subset \tilde{H}+\tilde{N}
$$

Similarly $\tilde{N}^{-}+\tilde{H}$ is a subalgebra of $\tilde{L}(A)$. We now consider the subspace $\tilde{N}^{-}+\tilde{H}+\tilde{N}$. The relations $\left[e_{i} f_{i}\right]=h_{i}$ and $\left[e_{i} f_{j}\right]=0$ if $i \neq j$ show that

$$
\left[e_{i}, \tilde{N}^{-}\right] \subset \tilde{N}^{-}+\tilde{H}
$$

For this is true for the generators of $\tilde{N}^{-}$, and the relation

$$
\left[e_{i}[x y]\right]=\left[\left[e_{i} x\right] y\right]+\left[x\left[e_{i} y\right]\right]
$$

then shows it is true for all elements of $\tilde{N}^{-}$since $\tilde{N}^{-}+\tilde{H}$ is a subalgebra. It follows that

$$
\left[e_{i}, \tilde{N}^{-}+\tilde{H}+\tilde{N}\right] \subset \tilde{N}^{-}+\tilde{H}+\tilde{N}
$$

since $\tilde{H}+\tilde{N}$ is a subalgebra. Similarly we have

$$
\left[f_{i}, \tilde{N}^{-}+\tilde{H}+\tilde{N}\right] \subset \tilde{N}^{-}+\tilde{H}+\tilde{N}
$$

and the relation

$$
\left[h_{i}, \tilde{N}^{-}+\tilde{H}+\tilde{N}\right] \subset \tilde{N}^{-}+\tilde{H}+\tilde{N}
$$

is clear. It follows that the set of all $x \in \tilde{L}(A)$ such that

$$
\left[x, \tilde{N}^{-}+\tilde{H}+\tilde{N}\right] \subset \tilde{N}^{-}+\tilde{H}+\tilde{N}
$$

contains $e_{i}, h_{i}, f_{i}$. However, the relation

$$
[[x y] z]=[[x z] y]+[x[y z]]
$$

for $z \in \tilde{N}^{-}+\tilde{H}+\tilde{N}$ shows that the set of such $x$ is a subalgebra. This subalgebra must be the whole of $\tilde{L}(A)$. Thus $\tilde{N}^{-}+\tilde{H}+\tilde{N}$ is an ideal of $\tilde{L}(A)$. Since $\tilde{L}(A)$ is generated by $e_{i}, h_{i}, f_{i}$ it follows that $\tilde{N}^{-}+\tilde{H}+\tilde{N}=\tilde{L}(A)$. We know that this sum is direct, so we have

$$
\tilde{L}(A)=\tilde{N}^{-} \oplus \tilde{H} \oplus \tilde{N}
$$

Since $\tilde{N}^{-} \subset \tilde{L}(A)^{-}, \tilde{N} \subset \tilde{L}(A)^{+}$and the sum $\tilde{L}(A)^{-}+\tilde{H}+\tilde{L}(A)^{+}$is direct we deduce that $\tilde{N}^{-}=\tilde{L}(A)^{-}$and $\tilde{N}=\tilde{L}(A)^{+}$. Since $\tilde{L}(A)=\tilde{L}(A)^{-} \oplus \tilde{H} \oplus$ $\tilde{L}(A)^{+}, \tilde{H} \subset \tilde{L}(A)_{0}$, and the weights occurring in $\tilde{L}(A)^{-}$and $\tilde{L}(A)^{+}$are all non-zero, we deduce from Proposition 7.11 that $\tilde{H}=\tilde{L}(A)_{0}$. Thus all parts of the proposition have been proved.

Proposition $7.13 \operatorname{dim} \tilde{L}(A)_{\alpha_{i}}=1$ and $\operatorname{dim} \tilde{L}(A)_{-\alpha_{i}}=1$
Proof. We know that $e_{i} \in \tilde{L}(A)_{\alpha_{i}}$. Also the element $e_{i} \in \tilde{L}(A)$ is non-zero, since it induces a non-zero endomorphism $\rho\left(e_{i}\right)$ on the $\tilde{L}(A)$-module $\mathfrak{F}^{-}$ considered in Proposition 7.9. Hence $\operatorname{dim} \tilde{L}(A)_{\alpha_{i}} \geq 1$. On the other hand we have

$$
\tilde{L}(A)_{\alpha_{i}} \subset \tilde{L}(A)^{+}=\tilde{N} .
$$

Now $\tilde{N}$ is generated by $e_{1}, \ldots, e_{l}$ so is spanned by monomials in these elements. All such monomials are weight vectors. The only monomial which has weight $\alpha_{i}$ is $e_{i}$, since the $\alpha_{i}$ are linearly independent. Thus we have $\operatorname{dim} \tilde{L}(A)_{\alpha_{i}}=1$. The relation $\operatorname{dim} \tilde{L}(A)_{-\alpha_{i}}=1$ is obtained similarly.

### 7.5 The existence theorem

We now turn to a study of the Lie algebra $L(A)$, in order to show that it is a finite dimensional simple Lie algebra with Cartan matrix $A$. From the definitions of $L(A), \tilde{L}(A)$ we see that $L(A)$ is isomorphic to $\tilde{L}(A) / I$ where $I$ is the ideal of $\tilde{L}(A)$ generated by the elements

$$
\begin{aligned}
& {\left[e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right]} \\
& {\left[f_{i}\left[f_{i} \ldots\left[f_{i} f_{j}\right]\right]\right]}
\end{aligned}
$$

for all $i \neq j$. As usual we have $1-A_{i j}$ factors $e_{i}$ or $f_{i}$.
Proposition 7.14 (i) Let $I^{+}$be the ideal of $\tilde{N}$ generated by the elements $\left[e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right]$ for all $i \neq j$. Then $I^{+}$is an ideal of $\tilde{L}(A)$.
(ii) Let $I^{-}$be the ideal of $\tilde{N}^{-}$generated by the elements $\left[f_{i}\left[f_{i} \ldots\left[f_{i} f_{j}\right]\right]\right]$ for all $i \neq j$. Then $I^{-}$is an ideal of $\tilde{L}(A)$.
(iii) $I=I^{+} \oplus I^{-}$.

Proof. We write $X_{i j}=\left[e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right]$ and $Y_{i j}=\left[f_{i}\left[f_{i} \ldots\left[f_{i} f_{j}\right]\right]\right]$. Then $I^{+}$ is the set of all linear combinations of elements

$$
\left[\left[X_{i j} e_{k_{1}}\right] \ldots e_{k_{r}}\right]
$$

for all $i \neq j$ and all $k_{1}, \ldots, k_{r}$ in $\{1, \ldots, l\}$. For such linear combinations certainly lie in $I^{+}$, and form an ideal of $\tilde{N}$.

Now $X_{i j}$ is a weight vector, being a Lie product of weight vectors $e_{i}, e_{j}$. Similarly $\left[\left[X_{i j} e_{k_{1}}\right] \ldots e_{k_{r}}\right]$ is a weight vector. It is therefore transformed by each of $h_{1}, \ldots, h_{l}$ into a scalar multiple of itself. In order to show that $I^{+}$is an ideal of $\tilde{L}(A)$ it will therefore be sufficient to show

$$
\left[f_{k},\left[\left[X_{i j} e_{k_{1}}\right] \ldots e_{k_{r}}\right]\right] \in I^{+}
$$

for all $i, j, k_{1}, \ldots, k_{r}, k$. We shall prove this by induction on $r$, beginning with $r=0$. In the following lemma we shall show that $\left[f_{k}, X_{i j}\right]=0$, thus beginning the induction. So let $r \geq 1$ and write $\left[\left[X_{i j} e_{k_{1}}\right] \ldots e_{k_{r-1}}\right]=y$. We assume $\left[f_{k} y\right] \in I^{+}$by induction. Then

$$
\left[f_{k}\left[y e_{k_{r}}\right]\right]=\left[\left[f_{k} y\right] e_{k_{r}}\right]+\left[y\left[f_{k} e_{k_{r}}\right]\right]
$$

If $k_{r} \neq k$ then $\left[f_{k} e_{k_{r}}\right]=0$ and so

$$
\left[f_{k}\left[y e_{k_{r}}\right]\right]=\left[\left[f_{k} y\right] e_{k_{r}}\right] \in I^{+}
$$

If $k_{r}=k$ then

$$
\left[f_{k}\left[y e_{k_{r}}\right]\right]=\left[\left[f_{k} y\right] e_{k_{r}}\right]+\left[h_{k} y\right] \in I^{+}
$$

This completes the induction. Thus $I^{+}$is an ideal of $\tilde{L}(A)$. Similarly $I^{-}$is an ideal of $\tilde{L}(A)$. Hence $I^{+} \oplus I^{-}$is an ideal of $\tilde{L}(A)$ containing the elements $X_{i j}$ and $Y_{i j}$. Moreover any ideal of $\tilde{L}(A)$ containing the $X_{i j}$ and $Y_{i j}$ must contain $I^{+}$and $I^{-}$. Hence $I^{+} \oplus I^{-}=I$.

In order to complete the proof of Proposition 7.14 we need the following lemma.

Lemma 7.15 $\left[f_{k}, X_{i j}\right]=0$ for all $i, j, k$ with $i \neq j$.
Proof. If $k \notin\{i, j\}$ this relation is obvious since $\left[f_{k} e_{i}\right]=0$ and $\left[f_{k} e_{j}\right]=0$.
So suppose $k=j$. Then we have

$$
\begin{aligned}
& {\left[f_{j}\left[e_{i} e_{j}\right]\right]=\left[e_{i}\left[f_{j} e_{j}\right]\right]=\left[h_{j} e_{i}\right]=A_{j i} e_{i}} \\
& {\left[f_{j}\left[e_{i}\left[e_{i} e_{j}\right]\right]\right]=\left[e_{i}\left[f_{j}\left[e_{i} e_{j}\right]\right]\right]=0} \\
& {\left[f_{j}\left[e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right]\right]=0 \quad \text { for } r \geq 2}
\end{aligned}
$$

by induction on $r$. Hence $\left[f_{j}, X_{i j}\right]=0$ if $1-A_{i j} \geq 2$, that is $A_{i j} \leq-1$. If $A_{i j}=0$ then $A_{j i}=0$ and $\left[f_{j}, X_{i j}\right]=0$ in this case also.

Finally suppose that $k=i$. In this case we shall show that, for $r \geq 1$,

$$
\left[f_{i}\left[e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right]\right]=-r\left(A_{i j}+r-1\right)\left[e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right] .
$$

For $r=1$ we have

$$
\left[f_{i}\left[e_{i} e_{j}\right]\right]=\left[\left[f_{i} e_{i}\right] e_{j}\right]=-\left[h_{i} e_{j}\right]=-A_{i j} e_{j}
$$

For $r>1$ we use induction. We have

$$
\begin{aligned}
& {\left[f_{i}\left[e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right]\right]} \\
& =-\left[h_{i}\left[e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right]\right]-(r-1)\left(A_{i j}+r-2\right)\left[e_{\leftarrow-1 \longrightarrow}^{e}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right] \\
& =\left(-\left(2 r-2+A_{i j}\right)-(r-1)\left(A_{i j}+r-2\right)\right)\left[\underset{\leftarrow r-1 \longrightarrow}{e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]}\right] \\
& =-r\left(A_{i j}+r-1\right)\left[\underset{\leftarrow r-1 \longrightarrow}{e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]}\right]
\end{aligned}
$$

as required. We now put $r=1-A_{i j}$ and obtain $\left[f_{i}, X_{i j}\right]=0$.
Corollary 7.16 $L(A)=N^{-} \oplus H \oplus N$ where $H$ is isomorphic to $\tilde{H}, N^{-}$is isomorphic to $\tilde{N}^{-} / I^{-}$and $N$ is isomorphic to $\tilde{N} / I^{+}$.

Proof. This follows from the facts that $L(A)$ is isomorphic to $\tilde{L}(A) / I, \tilde{L}(A)=$ $\tilde{N}^{-} \oplus \tilde{H} \oplus \tilde{N}$, and $I=I^{+} \oplus I^{-}$.

We shall continue to denote the generators of $L(A)$ by $e_{i}, h_{i}, f_{i}$. These are the images of the generators of $\mathfrak{R}$ under the natural homomorphism $\mathfrak{R} \rightarrow L(A)$.

Proposition 7.17 The maps ad $e_{i}: L(A) \rightarrow L(A)$ and ad $f_{i}: L(A) \rightarrow L(A)$ are locally nilpotent.

Proof. To show that ad $e_{i}$ is locally nilpotent we must show that, for all $x \in L(A)$, there exists $n(x)$ such that $\left(\operatorname{ad} e_{i}\right)^{n(x)} x=0$. Now if ad $e_{i}$ acts locally nilpotently on $x$ and $y$ it also acts locally nilpotently on $[x y]$. For

$$
\left(\operatorname{ad} e_{i}\right)^{n}[x y]=\sum_{r=0}^{n}\binom{n}{r}\left[\left(\operatorname{ad} e_{i}\right)^{r} x,\left(\operatorname{ad} e_{i}\right)^{n-r} y\right]
$$

$\left(\operatorname{ad} e_{i}\right)^{r} x$ will be 0 if $r$ is sufficiently large and $\left(\operatorname{ad} e_{i}\right)^{n-r} y$ will be 0 if $n-r$ is sufficiently large. Thus $\left(\operatorname{ad} e_{i}\right)^{n}[x y]$ will be 0 if $n$ is sufficiently large.

It follows that the set of elements of $L(A)$ on which ad $e_{i}$ acts locally nilpotently is a subalgebra. However, we have

$$
\begin{aligned}
& \operatorname{ad} e_{i} \cdot e_{i}=0 \\
& \left(\operatorname{ad} e_{i}\right)^{1-A_{i j}} e_{j}=0 \quad \text { if } i \neq j \\
& \left(\operatorname{ad} e_{i}\right)^{2} h_{j}=0 \quad \text { for all } j \\
& \left(\operatorname{ad} e_{i}\right)^{3} f_{i}=0 \\
& \operatorname{ad} e_{i} \cdot f_{j}=0 \quad \text { if } i \neq j
\end{aligned}
$$

Thus this subalgebra contains all the generators $e_{j}, h_{j}, f_{j}$ of $L(A)$, so is the whole of $L(A)$.

We see similarly that ad $f_{i}$ is locally nilpotent on $L(A)$.

Now the proof of Proposition 3.4 shows that if $\delta: L \rightarrow L$ is a locally nilpotent derivation of a Lie algebra $L$ then $\exp \delta$ is an automorphism of $L$. Thus expad $e_{i}$ and expad $f_{i}$ are automorphisms of $L(A)$. We define $\theta_{i} \in$ Aut $L(A)$ by

$$
\theta_{i}=\exp \operatorname{ad} e_{i} \cdot \exp \operatorname{ad}\left(-f_{i}\right) \cdot \exp \operatorname{ad} e_{i}
$$

Proposition 7.18 (i) $\theta_{i}(H)=H$
(ii) $\theta_{i}(h)=s_{i}(h)$ for all $h \in H$ where $s_{i}: H \rightarrow H$ is the linear map given by $s_{i}\left(h_{j}\right)=h_{j}-A_{j i} h_{i}$.

Proof. We have

$$
\begin{aligned}
& \exp \operatorname{ad} e_{i} \cdot h_{j}=\left(1+\operatorname{ad} e_{i}\right) h_{j}=h_{j}-A_{j i} e_{i} \\
& \exp \operatorname{ad}\left(-f_{i}\right) \cdot \exp \text { ad } e_{i} \cdot h_{j}=\exp \operatorname{ad}\left(-f_{i}\right) \cdot\left(h_{j}-A_{j i} e_{i}\right) \\
& =\left(1-\operatorname{ad} f_{i}+\frac{\left(\operatorname{ad} f_{i}\right)^{2}}{2}\right)\left(h_{j}-A_{j i} e_{i}\right) \\
& =h_{j}-A_{j i} e_{i}-A_{j i} f_{i}-A_{j i} h_{i}+A_{j i} f_{i}=h_{j}-A_{j i} h_{i}-A_{j i} e_{i} \\
& \exp \text { ad } e_{i} \cdot \exp \operatorname{ad}\left(-f_{i}\right) \cdot \exp \text { ad } e_{i} \cdot h_{j}=\exp \text { ad } e_{i}\left(h_{j}-A_{j i} h_{i}-A_{j i} e_{i}\right) \\
& =\left(1+\operatorname{ad} e_{i}\right)\left(h_{j}-A_{j i} h_{i}-A_{j i} e_{i}\right)=h_{j}-A_{j i} h_{i}-A_{j i} e_{i}-A_{j i} e_{i}+2 A_{j i} e_{i} \\
& =h_{j}-A_{j i} h_{i} .
\end{aligned}
$$

Now the action of $s_{i}$ on $H$ is precisely that of the fundamental reflection $s_{i}=s_{\alpha_{i}}$ defined in Section 6.4. We recall that

$$
\begin{aligned}
s_{i}(h) & =h-2 \frac{\left\langle h_{\alpha_{i}}^{\prime}, h\right\rangle}{\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle} h_{\alpha_{i}}^{\prime} \quad \text { for } h \in H \\
& =h-\left\langle h_{\alpha_{i}}^{\prime}, h\right\rangle h_{i}
\end{aligned}
$$

In particular

$$
s_{i}\left(h_{j}\right)=h_{j}-\left\langle h_{\alpha_{i}}^{\prime}, h_{j}\right\rangle h_{i}=h_{j}-2 \frac{\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle}{\left\langle h_{\alpha_{j}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle} h_{i}=h_{j}-A_{j i} h_{i}
$$

Thus Proposition 7.18 shows that the automorphism $\theta_{i}$ of $L(A)$ induces the fundamental reflection $s_{i}$ on $H$.

We now consider the decomposition of $L(A)$ into weight spaces with respect to $H$. This time the weights are elements of $\operatorname{Hom}(H, \mathbb{C})$. For each weight $\mu: H \rightarrow \mathbb{C}$ we define the weight space $L(A)_{\mu}$ by

$$
L(A)_{\mu}=\{x \in L(A) ;[h x]=\mu(h) x \quad \text { for all } h \in H\}
$$

Proposition 7.19 $L(A)=\bigoplus_{\mu} L(A)_{\mu}$.
Proof. The algebra $L(A)$ is the sum of its weight spaces, since its generators $e_{i}, h_{i}, f_{i}$ are weight vectors. Moreover the sum of weight spaces is direct, just as in the proof of Proposition 7.11.

It also follows from Proposition 7.12 and Corollary 7.16 that $L(A)=$ $N^{-} \oplus H \oplus N$ where all weights coming from $N$ are in $Q^{+}$and all weights coming from $N^{-}$are in $Q^{-}$. We also have $H=L(A)_{0}$.

Proposition $7.20 \operatorname{dim} L(A)_{\alpha_{i}}=1$ and $\operatorname{dim} L(A)_{-\alpha_{i}}=1$.
Proof. By Proposition 7.13 we certainly have $\operatorname{dim} L(A)_{\alpha_{i}} \leq 1$. However, the ideal $I^{+}$of $\tilde{N}$ such that $\tilde{N} / I^{+} \cong N$ has the property that $I^{+}$is a sum of weight spaces, and all weights occurring in $I^{+}$are sums of $\alpha_{1}, \ldots, \alpha_{l}$ involving at least two terms. This is clear from the proof of Proposition 7.14. Thus $\alpha_{i}$ is not a weight of $I^{+}$. Hence

$$
\operatorname{dim} L(A)_{\alpha_{i}}=\operatorname{dim} \tilde{L}(A)_{\alpha_{i}}=1
$$

One shows similarly that $\operatorname{dim} \tilde{L}(A)_{-\alpha_{i}}=1$.
Proposition 7.21 The automorphism $\theta_{i}$ of $L(A)$ transforms $L(A)_{\mu}$ to $L(A)_{s_{i} \mu}$. Hence $\operatorname{dim} L(A)_{\mu}=\operatorname{dim} L(A)_{s_{i} \mu}$.

Proof. Let $x \in L(A)_{\mu}$. Then $[h x]=\mu(h) x$ for all $h \in H$. We apply the automorphism $\theta_{i}$. This fixes $H$ by Proposition 7.18. We have

$$
\left[\theta_{i} h, \theta_{i} x\right]=\mu(h) \theta_{i} x
$$

Hence

$$
\left[h, \theta_{i} x\right]=\mu\left(\theta_{i}^{-1} h\right) \theta_{i} x=\mu\left(s_{i}^{-1} h\right) \theta_{i} x=\left(s_{i} \mu(h)\right) \theta_{i} x
$$

again by Proposition 7.18. Thus we have $\theta_{i} x \in L(A)_{s_{i} \mu}$. Hence

$$
\theta_{i}\left(L(A)_{\mu}\right) \subset L(A)_{s_{i} \mu}
$$

Replacing $\theta_{i}$ by $\theta_{i}^{-1}, \mu$ by $s_{i} \mu$ and recalling that $s_{i}^{2}=1$ we also obtain

$$
\theta_{i}^{-1}\left(L(A)_{s_{i} \mu}\right) \subset L(A)_{\mu}
$$

Hence $\theta_{i}\left(L(A)_{\mu}\right) \supset L(A)_{s_{i} \mu}$ and we have $\theta_{i}\left(L(A)_{\mu}\right)=L(A)_{s_{i} \mu}$.
We now define $W$ to be the group of non-singular linear transformations of $H^{*}=\operatorname{Hom}(H, \mathbb{C})$ generated by $s_{1}, \ldots, s_{l}$ and define $\Phi$ to be the set of elements $w\left(\alpha_{i}\right)$ for $w \in W$ and $i \in\{1, \ldots, l\}$. Then $\Phi$ is the root system determined by the given Cartan matrix $A$ and $W$ is its Weyl group.

Proposition 7.22 $\operatorname{dim} L(A)_{\alpha}=1$ for all $\alpha \in \Phi$.

Proof. We have $\alpha=w\left(\alpha_{i}\right)$ for some $i$ and some $w \in W$. Since $W$ is generated by $s_{1}, \ldots, s_{l}, w$ is a product of such elements. Thus it follows from Proposition 7.21 that

$$
\operatorname{dim} L(A)_{\alpha}=\operatorname{dim} L(A)_{\alpha_{i}}=1
$$

We aim to show that the Lie algebra $L(A)$ is finite dimensional. As a step in this direction we shall show that the Weyl group $W$ is finite. $W$ is isomorphic to the group of non-singular linear transformations of $H_{\mathbb{R}}$ generated by $s_{1}, \ldots, s_{l}$ where $H_{\mathbb{R}}=\mathbb{R} h_{1}+\cdots+\mathbb{R} h_{l}$. We have $\operatorname{dim} H_{\mathbb{R}}=l$. We do not have the scalar product on $H_{\mathbb{R}}$ available from the Killing form, so we define a scalar product directly from the Cartan matrix $A$.

Proposition 7.23 The Cartan matrix can be factorised as $A=D B$ where $D$ is diagonal and $B$ is symmetric. $D$ is the diagonal matrix with entries $d_{1}, \ldots, d_{l} \in\{1,2,3\}$ defined as follows.

If the Dynkin diagram has only single edges then all $d_{i}=1$.
If the Dynkin diagram has a double edge then $d_{i}=1$ if $\alpha_{i}$ is a long root and $d_{i}=2$ if $\alpha_{i}$ is a short root.

If the Dynkin diagram has a triple edge then $d_{i}=1$ if $\alpha_{i}$ is a long root and $d_{i}=3$ if $\alpha_{i}$ is a short root.

Proof. This may be checked from the standard list 6.12 of Cartan matrices.

For example in type $G_{2}$ we have

$$
\left(\begin{array}{cc}
2 & -1 \\
-3 & 2
\end{array}\right)=\left(\begin{array}{ll}
1 & 0 \\
0 & 3
\end{array}\right)\left(\begin{array}{cc}
2 & -1 \\
-1 & \frac{2}{3}
\end{array}\right)
$$

We now define a bilinear form on $H_{\mathbb{R}}$ by $\left\langle h_{i}, h_{j}\right\rangle=d_{i} d_{j} B_{i j}$. This form is symmetric since $B$ is a symmetric matrix.

Proposition 7.24 This scalar product is positive definite.
Proof. We have $n_{i j}=A_{i j} A_{j i}=d_{i} d_{j} B_{i j}^{2}$, thus $-\sqrt{ } n_{i j}=\sqrt{ } d_{i} \sqrt{ } d_{j} B_{i j}$. The matrix of our scalar product is

This matrix is congruent to the matrix $\left(\begin{array}{cccc}2 & & & \\ & \ddots \sqrt{ } n_{i j} \\ -\sqrt{ } n_{i j} & \cdot & \\ & & \end{array}\right)$ of the quadratic form $\mathrm{Q}\left(x_{1}, \ldots, x_{l}\right)$ of Proposition 6.6, which is positive definite. Thus $D B D$ is also positive definite.

Proposition 7.25 Our scalar product on $H_{\mathbb{R}}$ is invariant under $W$.

Proof. We first observe that

$$
\left\langle h_{i}, x\right\rangle=d_{i} \alpha_{i}(x) \quad \text { for all } x \in H_{\mathbb{R}}
$$

For $\left\langle h_{i}, h_{j}\right\rangle=d_{i} d_{j} B_{i j}=d_{i} A_{j i}=d_{i} \alpha_{i}\left(h_{j}\right)$. It is sufficient to show that $\left\langle s_{i} x, s_{i} y\right\rangle=\langle x, y\rangle$ for all $x, y \in H_{\mathbb{R}}$. We note that $s_{i}(x)=x-\alpha_{i}(x) h_{i}$ since $s_{i}\left(h_{j}\right)=h_{j}-A_{j i} h_{i}$. Hence

$$
\begin{aligned}
\left\langle s_{i} x, s_{i} y\right\rangle & =\left\langle x-\alpha_{i}(x) h_{i}, y-\alpha_{i}(y) h_{i}\right\rangle \\
& =\langle x, y\rangle-\alpha_{i}(x)\left\langle h_{i}, y\right\rangle-\alpha_{i}(y)\left\langle h_{i}, x\right\rangle+\alpha_{i}(x) \alpha_{i}(y)\left\langle h_{i}, h_{i}\right\rangle \\
& =\langle x, y\rangle-d_{i} \alpha_{i}(x) \alpha_{i}(y)-d_{i} \alpha_{i}(x) \alpha_{i}(y)+2 d_{i} \alpha_{i}(x) \alpha_{i}(y)=\langle x, y\rangle
\end{aligned}
$$

Thus the Weyl group $W$ acts as a group of isometries on the Euclidean space $H_{\mathbb{R}}$.

We define certain subsets of $H_{\mathbb{R}}$ as follows:

$$
\begin{aligned}
& H_{i}=\left\{x \in H_{\mathbb{R}} ;\left\langle h_{i}, x\right\rangle=0\right\} \\
& H_{i}^{+}=\left\{x \in H_{\mathbb{R}} ;\left\langle h_{i}, x\right\rangle>0\right\} \\
& H_{i}^{-}=\left\{x \in H_{\mathbb{R}} ;\left\langle h_{i}, x\right\rangle<0\right\} \\
& C=H_{1}^{+} \cap \cdots \cap H_{l}^{+} .
\end{aligned}
$$

$C$ is called the fundamental chamber.
Let $W_{i j}$ be the subgroup of $W$ generated by $s_{i}, s_{j}$, where $i \neq j . s_{i} s_{j}$ has finite order $m_{i j}$ given in terms of the Cartan matrix by $2 \cos \left(\pi / m_{i j}\right)=\sqrt{ } n_{i j}$. Thus $W_{i j}$ is a finite dihedral group.

Lemma 7.26 Let $w \in W_{i j}$ with $i \neq j$. Then either $w\left(H_{i}^{+} \cap H_{j}^{+}\right) \subset H_{i}^{+}$or $w\left(H_{i}^{+} \cap H_{j}^{+}\right) \subset H_{i}^{-}$and $l\left(s_{i} w\right)=l(w)-1$.

Proof. Let $U$ be the 2-dimensional subspace of $H_{\mathbb{R}}$ spanned by $h_{i}, h_{j}$ and $U^{\perp}$ be the orthogonal subspace. Then $H_{\mathbb{R}}=U \oplus U^{\perp}$ and the elements of $W_{i j}$ act trivially on $U^{\perp}$. It is therefore sufficient to prove the result in $U$. Let $\Gamma=U \cap H_{i}^{+} \cap H_{j}^{+}$. We obtain a configuration of chambers in $U$ as shown in Figure 7.1.


Figure 7.1 Configuration of chambers in $U$

The chambers $\Gamma, s_{j}(\Gamma), s_{j} s_{i}(\Gamma), \ldots, s_{j} s_{i} \ldots(\Gamma)$ lie in $H_{i}^{+}$and the cham$\leftarrow m_{i j}-1 \longrightarrow$ bers

$$
s_{i}(\Gamma), s_{i} s_{j}(\Gamma), s_{i} s_{j} s_{i}(\Gamma), \ldots, \begin{aligned}
& s_{i} s_{j} \ldots(\Gamma) \\
& \leftarrow m_{i j} \longrightarrow
\end{aligned}
$$

lie in $H_{i}^{-}$. The elements $s_{i}, s_{i} s_{j}, \ldots, s_{i} s_{j} \ldots$ of $W_{i j}$ all satisfy $l\left(s_{i} w\right)=l(w)-$ $\longleftarrow m_{i j} \longrightarrow$

1. Thus for each $w \in W_{i j}$ we have either $w\left(H_{i}^{+} \cap H_{j}^{+}\right) \subset H_{i}^{+}$or $w\left(H_{i}^{+} \cap H_{j}^{+}\right) \subset$ $H_{i}^{-}$and $l\left(s_{i} w\right)=l(w)-1$.

Proposition 7.27 (a) Let $w \in W$. Then either $w(C) \subset H_{i}^{+}$or $w(C) \subset H_{i}^{-}$and $l\left(s_{i} w\right)=l(w)-1$.
(b) Let $w \in W$ and $i \neq j$. Then there exists $w^{\prime} \in W_{i j}$ such that $w(C) \subset$ $w^{\prime}\left(H_{i}^{+} \cap H_{j}^{+}\right)$and $l(w)=l\left(w^{\prime}\right)+l\left(w^{\prime-1} w\right)$.
Note. Part (a) is the result we shall need. To prove it we must also prove part (b) at the same time.

Proof. We prove both statements together by induction on $l(w)$. If $l(w)=0$ then $w=1$ and (a), (b) are true. So suppose $l(w)>0$. Then $w=s_{j} w^{\prime}$ with $l\left(w^{\prime}\right)=l(w)-1$ for some $j$. We prove (a).

First suppose $j=i$. By induction $w^{\prime}(C) \subset H_{i}^{+}$or $w^{\prime}(C) \subset H_{i}^{-}$and $l\left(s_{i} w^{\prime}\right)=$ $l\left(w^{\prime}\right)-1$. But $l\left(s_{i} w^{\prime}\right)=l\left(w^{\prime}\right)+1$, so $w^{\prime}(C) \subset H_{i}^{+}$. Then $w(C) \subset H_{i}^{-}$and $l\left(s_{i} w\right)=l(w)-1$.

Now suppose $j \neq i$. By induction there exists $w^{\prime \prime} \in W_{i j}$ with $w^{\prime}(C) \subset$ $w^{\prime \prime}\left(H_{i}^{+} \cap H_{j}^{+}\right)$and $l\left(w^{\prime}\right)=l\left(w^{\prime \prime}\right)+l\left(w^{\prime \prime-1} w^{\prime}\right)$. Thus $w(C) \subset s_{j} w^{\prime \prime}\left(H_{i}^{+} \cap H_{j}^{+}\right)$. By Lemma 7.26 we have either $s_{j} w^{\prime \prime}\left(H_{i}^{+} \cap H_{j}^{+}\right) \subset H_{i}^{+}$or $s_{j} w^{\prime \prime}\left(H_{i}^{+} \cap H_{j}^{+}\right) \subset$ $H_{i}^{-}$and $l\left(s_{i} s_{j} w^{\prime \prime}\right)=l\left(s_{j} w^{\prime \prime}\right)-1$. In the first case $w(C) \subset H_{i}^{+}$. In the second case $\quad w(C) \subset H_{i}^{-} \quad$ and $\quad l\left(s_{i} w\right)=l\left(s_{i} s_{j} w^{\prime}\right)=l\left(s_{i} s_{j} w^{\prime \prime} w^{\prime \prime-1} w^{\prime}\right) \leq l\left(s_{i} s_{j} w^{\prime \prime}\right)+$ $l\left(w^{\prime \prime-1} w^{\prime}\right)=l\left(s_{j} w^{\prime \prime}\right)-1+l\left(w^{\prime \prime-1} w^{\prime}\right) \leq l\left(w^{\prime \prime}\right)+l\left(w^{\prime \prime-1} w^{\prime}\right)=l\left(w^{\prime}\right)=l(w)-1$. Thus $l\left(s_{i} w\right)=l(w)-1$ and (a) is proved.

We now prove (b). If $w(C) \subset H_{i}^{+} \cap H_{j}^{+}$then (b) holds with $w^{\prime}=1$. Thus we may assume without loss of generality that $w(C) \not \subset H_{i}^{+}$. So by (a), which is now proved under the assumption of the inductive hypothesis, $w(C) \subset H_{i}^{-}$ and $l\left(s_{i} w\right)=l(w)-1$. By induction there exists $w^{\prime} \in W_{i j}$ such that $s_{i} w(C) \subset$ $w^{\prime}\left(H_{i}^{+} \cap H_{j}^{+}\right)$and $l\left(s_{i} w\right)=l\left(w^{\prime}\right)+l\left(w^{\prime-1} s_{i} w\right)$. Thus $w(C) \subset s_{i} w^{\prime}\left(H_{i}^{+} \cap H_{j}^{+}\right)$ and

$$
l(w)=1+l\left(s_{i} w\right)=1+l\left(w^{\prime}\right)+l\left(w^{\prime-1} s_{i} w\right) \geq l\left(s_{i} w^{\prime}\right)+l\left(\left(s_{i} w^{\prime}\right)^{-1} w\right) \geq l(w)
$$

Thus we have equality throughout and $l(w)=l\left(s_{i} w^{\prime}\right)+l\left(\left(s_{i} w^{\prime}\right)^{-1} w\right)$. Hence $s_{i} w^{\prime} \in W_{i j}$ is the required element and (b) is proved.

Proposition 7.28 If $w \in W$ satisfies $C \cap w(C)$ is non-empty, then $w=1$.

Proof. Suppose $w \neq 1$. Then $w=s_{i} w^{\prime}$ with $l\left(w^{\prime}\right)=l(w)-1$. By Proposition 7.27 (a) $w^{\prime}(C) \subset H_{i}^{+}$. Thus $w(C) \subset H_{i}^{-}$. Hence

$$
C \cap w(C) \subset H_{i}^{+} \cap H_{i}^{-}=\varnothing
$$

So if $C \cap w(C)$ is non-empty, $w=1$.
Now the Euclidean space $H_{\mathbb{R}}$ has an orthonormal basis, and the isometries of $H_{\mathbb{R}}$ are represented by orthogonal matrices with respect to this basis. Thus $W \subset O_{l}$ where $O_{l}$ is the group of $l \times l$ orthogonal matrices. $O_{l} \subset M_{l}$, the set of all $l \times l$ matrices over $\mathbb{R}$.

For any matrix $M=\left(m_{i j}\right) \in M_{l}$ we define $\|M\|=\sqrt{\sum_{i, j} m_{i j}^{2}}$ and for any column vector $v=\left(\lambda_{1}, \ldots, \lambda_{l}\right)^{\mathrm{t}} \in \mathbb{R}^{l}$ we define $\|v\|=\sqrt{\sum_{i} \lambda_{i}^{2}}$.

Lemma 7.29 (a) If $M \in O_{l}, v \in \mathbb{R}^{l}$ then $\|M v\|=\|v\|$.
(b) If $M \in O_{l}, N \in M_{l}$ then $\|M N\|=\|N\|$.
(c) If $M \in M_{l}, v \in \mathbb{R}^{l}$ then $\|M v\| \leq\|M\|\|v\|$.

Proof. Straightforward.

Proposition 7.30 (a) $W$ is finite.
(b) $\Phi$ is finite.

Proof. Since $\Phi=W(\Pi)$ where $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$ it is clear that (a) implies (b).
Thus we show that $W$ is finite. We consider the $W$-action on the Euclidean space $H_{\mathbb{R}}$. We give elements of $H_{\mathbb{R}}$ coordinates relative to our orthonormal basis. Let $v=\left(\lambda_{1}, \ldots, \lambda_{l}\right)^{t} \in C$. By definition of $C$ there exists $r>0$ such that $B_{r}(v) \subset C$ where

$$
B_{r}(v)=\left\{x \in \mathbb{R}^{l} ;\|x-v\|<r\right\} .
$$

Let $w \in W$ with $w \neq 1$. Then $w(C) \cap C=\varnothing$ by Proposition 7.28. Thus $w(v) \notin C$ so

$$
\|w(v)-v\| \geq r
$$

Hence

$$
\|w-1\|\|v\| \geq\|w(v)-v\| \geq r
$$

so $\|w-1\| \geq \frac{r}{\|v\|}$. Put $\varepsilon=\frac{r}{\|v\|}$. Then $\|w-1\| \geq \varepsilon$ for all $w \neq 1$ in $W$. Now let $w, w^{\prime} \in W$ have $w \neq w^{\prime}$. Then

$$
\left\|w-w^{\prime}\right\|=\left\|w^{\prime}\left(w^{\prime-1} w-1\right)\right\|=\left\|w^{\prime-1} w-1\right\| \geq \varepsilon
$$

since $w^{\prime} \in O_{l}$. Thus distinct elements of $W$ are separated by a distance of at least $\varepsilon$. Since $O_{l}$, and hence $W$, is bounded it follows that $W$ is finite.

We now return to our Lie algebra $L(A)$. We know that $\operatorname{dim} L(A)_{0}=l$ and $\operatorname{dim} L(A)_{\alpha}=1$ for all $\alpha \in \Phi$.

If we can prove that $L(A)_{\mu}=O$ for all $\mu \in H^{*}$ with $\mu \notin \Phi \cup\{0\}$ we shall be able to deduce that $L(A)$ is finite dimensional.

Proposition 7.31 Suppose $\mu \in H^{*}$ satisfies $\mu \neq 0$ and $\mu \notin \Phi$. Then $L(A)_{\mu}=O$.

Proof. We assume that $\mu \neq 0$ and $L(A)_{\mu} \neq O$. Since $\operatorname{dim} L(A)_{\mu} \leq \operatorname{dim} \tilde{L}(A)_{\mu}$ we see by Proposition 7.12 (iii) that $\mu \in Q^{+}$or $\mu \in Q^{-}$. In particular $\mu$ lies in the vector space $H_{\mathbb{R}}^{*}$ of real linear combinations of $\alpha_{1}, \ldots, \alpha_{l}$.

Suppose first that $\mu$ is a multiple of some root $\alpha \in \Phi$. Then $\mu=n \alpha$ or $-n \alpha$ with $n>0$ and $\alpha \in \Phi^{+}$. Now $\alpha=w\left(\alpha_{i}\right)$ for some $w \in W$ and some $\alpha_{i} \in \Pi$ and we have $\operatorname{dim} L(A)_{n \alpha}=\operatorname{dim} L(A)_{n \alpha_{i}}$ by Proposition 7.21. Hence $L(A)_{n \alpha_{i}} \neq O$. Now $N$ is generated by elements $e_{1}, \ldots, e_{l}$ and no non-zero Lie product of these can have weight $n \alpha_{i}$ unless $n=1$. Thus $\mu=\alpha$ or $-\alpha$, that is $\mu \in \Phi$.

Secondly suppose $\mu$ is not a multiple of a root. Let

$$
\begin{aligned}
& H_{\mu}=\left\{h \in H_{\mathbb{R}} ; \mu(h)=0\right\} \\
& H_{\alpha}=\left\{h \in H_{\mathbb{R}} ; \alpha(h)=0\right\} .
\end{aligned}
$$

Then $H_{\mu}$ is distinct from all the $H_{\alpha}, \alpha \in \Phi$. Since $\Phi$ is finite we can find $h \in H_{\mu}$ such that $h \notin H_{\alpha}$ for all $\alpha \in \Phi$. It follows that $w(h) \notin H_{\alpha}$ for all $\alpha \in \Phi$, since $W$ permutes the $H_{\alpha}$.

We claim there exists $w \in W$ such that $\alpha_{i}(w(h))>0$ for all $i=1, \ldots, l$. To see this we define the height of an element of $H_{\mathbb{R}}$ by

$$
\operatorname{ht}\left(\sum n_{i} h_{i}\right)=\sum n_{i} .
$$

We choose an element $w \in W$ such that ht $w(h)$ is maximum. This is certainly possible as $W$ is finite. Then

$$
s_{i}(w(h))=w(h)-\alpha_{i}(w(h)) h_{i}
$$

Since ht $s_{i}(w(h)) \leq$ ht $w(h)$ we have $\alpha_{i}(w(h)) \geq 0$. However, $\alpha_{i}(w(h))=0$ would imply $w(h) \in H_{\alpha_{i}}$ which is impossible. Thus $\alpha_{i}(w(h))>0$ for each $i$ and so $w(h) \in C$.

Now we have $(w(\mu))(w(h))=\mu(h)=0$. We write $w(\mu)=\sum_{i=1}^{l} m_{i} \alpha_{i}$. Then we have

$$
\sum_{i=1}^{l} m_{i} \alpha_{i}(w(h))=0
$$

Since $\alpha_{i}(w(h))>0$ for each $i$ we must have some $m_{i}>0$ and some $m_{j}<0$ in the sum. Thus $w(\mu) \notin Q^{+}$and $w(\mu) \notin Q^{-}$. Hence $\tilde{L}(A)_{w(\mu)}=O$. By Proposition 7.21 we deduce $\tilde{L}(A)_{\mu}=O$, and so $L(A)_{\mu}=O$. This gives the required contradiction.

Corollary 7.32 (i) $L(A)=H \oplus \sum_{\alpha \in \Phi} L(A)_{\alpha}$.
(ii) $\operatorname{dim} L(A)=l+|\Phi|$.

Proof. This is evident since $L(A)$ is the direct sum of its weight spaces. The 0 -weight space is $H$ and this has dimension $l$. The only non-zero weights are the elements of $\Phi$ and the corresponding weight spaces are 1-dimensional. Thus we have the required formula for the dimension of $L(A)$.

We now know that $L(A)$ is a finite dimensional Lie algebra - indeed it has the dimension required for a simple Lie algebra with Cartan matrix $A$. We shall now be readily able to show that $L(A)$ has the required properties.

Proposition 7.33 The Lie algebra $L(A)$ is semisimple.

Proof. Let $R$ be the soluble radical of $L(A)$ and consider the series

$$
R=R^{(0)} \supset R^{(1)} \supset \cdots \supset R^{(n-1)} \supset R^{(n)}=O
$$

where $R^{(i+1)}=\left[R^{(i)} R^{(i)}\right]$. We write $I=R^{(n-1)}$. We suppose if possible that $R \neq O$. Then $I$ is a non-zero abelian ideal of $L$. Moreover $I$ is invariant under all automorphisms of $L$.

Since $[H I] \subset I$ we may regard $I$ as an $H$-module. We decompose it into its weight spaces. This weight space decomposition is

$$
I=(H \cap I) \oplus \sum_{\alpha \in \Phi}\left(L_{\alpha} \cap I\right) .
$$

For let $x \in I$ have $x=x_{0}+\sum_{\alpha \in \Phi} x_{\alpha}$ where $x_{0} \in H$ and $x_{\alpha} \in L_{\alpha}$. We show $x_{0} \in I$ and each $x_{\alpha} \in I$. There exists $h \in H$ such that $\alpha(h) \neq 0$ and $\beta(h) \neq \alpha(h)$ for all $\beta \in \Phi$ with $\beta \neq \alpha$. Then

$$
\operatorname{ad} \mathrm{h} \prod_{\substack{\beta \in \Phi \\ \beta \neq \alpha}}(\operatorname{ad~h}-\beta(h) 1) x=\alpha(h) \prod_{\substack{\beta \in \Phi \\ \beta \neq \alpha}}(\alpha(h)-\beta(h)) x_{\alpha}
$$

and this is an element of $I$. Hence $x_{\alpha} \in I$. Since this is true for each $\alpha \in \Phi$ we also have $x_{0} \in I$. Hence

$$
I=(H \cap I) \oplus \sum_{\alpha \in \Phi}\left(L_{\alpha} \cap I\right) .
$$

We claim that $L_{\alpha} \cap I=O$ for each $\alpha \in \Phi$. Otherwise we would have $L_{\alpha} \subset I$. Now $\alpha=w\left(\alpha_{i}\right)$ for some $w \in W$ and some $i=1, \ldots, l$. By Proposition 7.21 we can find an automorphism of $L(A)$ which transforms $L_{\alpha}$ to $L_{\alpha_{i}}$. Since $I$ is invariant under all automorphism we would obtain $L_{\alpha_{i}} \subset I$. Hence $e_{i} \in I$. But then $\left[e_{i} f_{i}\right]=h_{i} \in I$ and we would have $\left[h_{i} e_{i}\right]=2 e_{i} \in I$, contradicting the fact that $I$ is abelian. Hence $L_{\alpha} \cap I=O$ for all $\alpha \in \Phi$ and so $I \subset H$. Let $x \in I$. Then $\left[x e_{i}\right]=\alpha_{i}(x) e_{i} \in I$ hence $\alpha_{i}(x)=0$. Since $\alpha_{1}, \ldots, \alpha_{l}$ are linearly independent we deduce that $x=0$. Hence $I=O$, a contradiction.

Proposition 7.34 $H$ is a Cartan subalgebra of $L(A)$.

Proof. Since $H$ is abelian it is sufficient to show that $H=N(H)$. Let $x \in N(H)$. Then $x=h^{\prime}+\sum_{\alpha \in \Phi} \lambda_{\alpha} e_{\alpha}$ for $h^{\prime} \in H, e_{\alpha} \in L_{\alpha}$. Then for all $h \in H$ we have

$$
[h x]=\sum_{\alpha \in \Phi} \lambda_{\alpha} \alpha(h) e_{\alpha} \in H
$$

However, we can find $h \in H$ such that $\alpha(h) \neq 0$ for all $\alpha \in \Phi$. We deduce that $\lambda_{\alpha}=0$ for all $\alpha \in \Phi$, hence $x \in H$.

Proposition $7.35 L(A)$ is a simple Lie algebra with Cartan matrix A.
Proof. $L(A)=H \oplus \sum_{\alpha \in \Phi} L(A)_{\alpha}$ is the Cartan decomposition of $L(A)$ with respect to $H$. Thus $\Phi$ is the root system of $L(A)$. The Cartan matrix $A^{\prime}=\left(A_{i j}^{\prime}\right)$ of $L(A)$ is determined by the condition

$$
s_{i}\left(\alpha_{j}\right)=\alpha_{j}-A_{i j}^{\prime} \alpha_{i} .
$$

However, we have

$$
s_{i}\left(h_{j}\right)=h_{j}-A_{j i} h_{i} \quad \text { by Proposition 7.18. }
$$

Using the facts that $\left(s_{i} \alpha_{j}\right) h_{k}=\alpha_{j}\left(s_{i} h_{k}\right)$ and $\alpha_{j}\left(h_{k}\right)=A_{k j}$ we deduce that

$$
s_{i}\left(\alpha_{j}\right)=\alpha_{j}-A_{i j} \alpha_{i} .
$$

Hence $A^{\prime}=A$ and the Cartan matrix of $L(A)$ is $A$.
Since the Dynkin diagram of $A$ is assumed connected, $L(A)$ must be a simple Lie algebra, by Proposition 6.13.

Thus we have constructed, for each Cartan matrix on the standard list 6.12 a finite dimensional simple Lie algebra $L(A)$ with Cartan matrix $A$.

Theorem 7.36 The finite dimensional non-trivial simple Lie algebras over $\mathbb{C}$ are

| $A_{l}$ | $l \geq 1$ |
| :--- | :--- |
| $B_{l}$ | $l \geq 2$ |
| $C_{l}$ | $l \geq 3$ |
| $D_{l}$ | $l \geq 4$ |
| $E_{6}, E_{7}, E_{8}$ |  |
| $F_{4}$ |  |
| $G_{2}$ |  |

These Lie algebras are pairwise non-isomorphic.
Proof. For each Cartan matrix on the standard list 6.12 there is a corresponding finite dimensional simple Lie algebra, which by Theorem 7.5 is determined up to isomorphism. Simple Lie algebras with different Cartan matrices cannot be isomorphic since, by Proposition 6.4, the Cartan matrix on the standard list is uniquely determined by the Lie algebra.

The description in Proposition 7.35 of the simple Lie algebras by generators and relations enables us to choose the root vectors $e_{\alpha}$ in a way which makes the structure constants $N_{\alpha, \beta}$ very simple.

Theorem 7.37 It is possible to choose the root vectors $e_{\alpha}$ in the simple Lie algebra $L(A)$ in such a way that $N_{\alpha, \beta}= \pm(p+1)$ where $-p \alpha+\beta, \ldots, q \alpha+\beta$ are the $\alpha$-chain of roots through $\beta$.

Proof. $L(A)$ is the Lie algebra generated by elements $e_{1}, \ldots, e_{l}$, $h_{1}, \ldots, h_{l}, f_{1}, \ldots f_{l}$ subject to relations

$$
\begin{aligned}
& {\left[h_{i} h_{j}\right]=0} \\
& {\left[h_{i} e_{j}\right]=A_{i j} e_{j}} \\
& {\left[h_{i} f_{j}\right]=-A_{i j} f_{j}} \\
& {\left[e_{i} f_{i}\right]=h_{i}} \\
& {\left[e_{i} f_{j}\right]=0 \quad \text { if } i \neq j} \\
& {\left[e_{i}\left[e_{i} \ldots\left[e_{i} e_{j}\right]\right]\right]=0 \quad \text { if } i \neq j} \\
& {\left[f_{i}\left[f_{i} \ldots\left[f_{i} f_{j}\right]\right]\right]=0 \quad \text { if } i \neq j}
\end{aligned}
$$

with $1-A_{i j}$ occurrences of $e_{i}, f_{i}$ respectively.
We now define

$$
e_{i}^{\prime}=-f_{i}, \quad h_{i}^{\prime}=-h_{i}, \quad f_{i}^{\prime}=-e_{i}
$$

It is straightforward to check that $e_{i}^{\prime}, h_{i}^{\prime}, f_{i}^{\prime}$ satisfy the above relations. Thus there is a homomorphism $\omega: L(A) \rightarrow L(A)$ satisfying $\omega\left(e_{i}\right)=-f_{i}$, $\omega\left(h_{i}\right)=-h_{i}, \omega\left(f_{i}\right)=-e_{i}$. Since $\omega^{2}=1, \omega$ is an automorphism of $L(A)$.

Let $\alpha$ be a positive root of $L(A)$. Then

$$
\left[h_{i} e_{\alpha}\right]=\alpha\left(h_{i}\right) e_{\alpha}
$$

and so

$$
\left[-h_{i}, \theta\left(e_{\alpha}\right)\right]=\alpha\left(h_{i}\right) \theta\left(e_{\alpha}\right)
$$

that is

$$
\left[h_{i}, \theta\left(e_{\alpha}\right)\right]=-\alpha\left(h_{i}\right) \theta\left(e_{\alpha}\right)
$$

whence $\theta\left(e_{\alpha}\right) \in L_{-\alpha}$. Let $\theta\left(e_{\alpha}\right)=\lambda e_{-\alpha}$. Then $\lambda \neq 0$ and we may choose $\mu \in \mathbb{C}$ with $\mu^{2}=-\lambda^{-1}$. Then

$$
\theta\left(\mu e_{\alpha}\right)=-\mu^{-1} e_{-\alpha}
$$

and

$$
\left[\mu e_{\alpha}, \mu^{-1} e_{-\alpha}\right]=\left[e_{\alpha} e_{-\alpha}\right]=\frac{2 h_{\alpha}^{\prime}}{\left\langle h_{\alpha}^{\prime}, h_{\alpha}^{\prime}\right\rangle}
$$

We now change our choice of the root vectors $e_{\alpha} \in L_{\alpha}$. For each positive root $\alpha$ we take $\mu e_{\alpha}$ as our root vector and for the corresponding negative root $-\alpha$ we take $\mu^{-1} e_{-\alpha}$ as the root vector. Changing the notation to call these
new root vectors $e_{\alpha}, e_{-\alpha}$ we retain the multiplication formulae of Section 7.1, except that the structure constants $N_{\alpha, \beta}$ may now be altered. We also now have $\omega\left(e_{\alpha}\right)=-e_{-\alpha}$. Now

$$
\left[e_{\alpha} e_{\beta}\right]=N_{\alpha, \beta} e_{\alpha+\beta}
$$

and so

$$
\left[-e_{-\alpha},-e_{-\beta}\right]=N_{\alpha, \beta}\left(-e_{-\alpha-\beta}\right)
$$

This implies $N_{-\alpha,-\beta}=-N_{\alpha, \beta}$. By Proposition 7.1 (iii) we have $N_{\alpha, \beta} N_{-\alpha,-\beta}=$ $-(p+1)^{2}$, where $-p \alpha+\beta, \ldots, q \alpha+\beta$ is the $\alpha$-chain of roots through $\beta$. Hence $N_{\alpha, \beta}^{2}=(p+1)^{2}$ and $N_{\alpha, \beta}= \pm(p+1)$.

This result has important implications in the theory of Chevalley groups over arbitrary fields. (See, for example, R. W. Carter, Simple Groups of Lie Type, Wiley Classics Library, 1989.)

The signs in the formula $N_{\alpha, \beta}= \pm(p+1)$ can be chosen in various ways. By Lemma 7.3 and Proposition 7.4 the signs can be chosen arbitrarily for extraspecial pairs of roots $(\alpha, \beta)$ and are then determined for all other pairs $(\alpha, \beta)$.

## 8

## The simple Lie algebras

Having obtained a classification of the finite dimensional simple Lie algebras over $\mathbb{C}$ we shall in the present chapter investigate them individually in order to obtain their dimensions and a description of their root systems. In the case of Lie algebras of type $A_{l}, B_{l}, C_{l}$ or $D_{l}$ we shall also give a description in terms of Lie algebras of matrices.

The strategy for obtaining this information will be as follows. Given a Cartan matrix $A$ on the standard list 6.12 we shall describe a symmetric scalar product $\{$,$\} on an l$-dimensional vector space $V$ over $\mathbb{R}$ with basis $\alpha_{1}, \ldots, \alpha_{l}$ such that

$$
2 \frac{\left\{\alpha_{i}, \alpha_{j}\right\}}{\left\{\alpha_{i}, \alpha_{i}\right\}}=A_{i j} \quad i, j=1, \ldots, l
$$

We compare this scalar product with the Killing form $\left\langle\alpha_{i}, \alpha_{j}\right\rangle$ obtained when $\alpha_{1}, \ldots, \alpha_{l}$ are interpreted as a fundamental system of roots in the simple Lie algebra with Cartan matrix $A$. We claim there exists a constant $\kappa$ such that

$$
\left\langle\alpha_{i}, \alpha_{j}\right\rangle=\kappa\left\{\alpha_{i}, \alpha_{j}\right\} \quad \text { for all } i, j
$$

In fact we can define $\kappa$ by the equation

$$
\left\langle\alpha_{1}, \alpha_{1}\right\rangle=\kappa\left\{\alpha_{1}, \alpha_{1}\right\}
$$

Then we have

$$
2 \frac{\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}=2 \frac{\left\{\alpha_{i}, \alpha_{j}\right\}}{\left\{\alpha_{i}, \alpha_{i}\right\}} \quad \text { for all } i, j
$$

and since both scalar products are symmetric we deduce that

$$
\frac{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}=\frac{\left\{\alpha_{j}, \alpha_{j}\right\}}{\left\{\alpha_{i}, \alpha_{i}\right\}} \quad \text { for all } i, j
$$

Putting $j=1$ we deduce

$$
\left\langle\alpha_{i}, \alpha_{i}\right\rangle=\kappa\left\{\alpha_{i}, \alpha_{i}\right\} \quad \text { for all } i
$$

and it follows that

$$
\left\langle\alpha_{i}, \alpha_{j}\right\rangle=\kappa\left\{\alpha_{i}, \alpha_{j}\right\} \quad \text { for all } i, j
$$

Thus the symmetric scalar product $\{$,$\} is the same as the Killing form up$ to multiplication by the constant $\kappa$. In practise it will not be necessary to determine this constant.

We then consider the fundamental reflections $s_{i}: V \rightarrow V$ defined by

$$
s_{i}\left(\alpha_{j}\right)=\alpha_{j}-A_{i j} \alpha_{i}
$$

The maps $s_{1}, \ldots, s_{l}$ generate the Weyl group $W$ of transformations of $V$. The vectors in $V$ of form $w\left(\alpha_{i}\right)$ for all $w \in W$ and all $i$ will then give the full root system $\Phi$. We shall then be able to obtain the dimension of the simple Lie algebra $L$ by the formula

$$
\operatorname{dim} L=l+|\Phi|
$$

### 8.1 Lie algebras of type $A_{l}$

It will be convenient to describe the vector space $V$ as a subspace of a larger vector space $\tilde{V}$ of dimension $l+1$.

Let $\tilde{V}$ be a vector space over $\mathbb{R}$ with basis $\beta_{1}, \ldots, \beta_{l+1}$ and let the symmetric scalar product $\{$,$\} on \tilde{V}$ be defined by

$$
\left\{\beta_{i}, \beta_{j}\right\}=\delta_{i j} \quad i, j=1, \ldots, l+1
$$

We define $\alpha_{1}, \ldots, \alpha_{l}$ by

$$
\alpha_{1}=\beta_{1}-\beta_{2}, \quad \alpha_{2}=\beta_{2}-\beta_{3}, \quad \ldots, \quad \alpha_{l}=\beta_{l}-\beta_{l+1}
$$

Let $V$ be the subspace of $\tilde{V}$ spanned by $\alpha_{1}, \ldots, \alpha_{l}$. Then we have $\operatorname{dim} V=l$. Our scalar product satisfies

$$
\left\{\alpha_{i}, \alpha_{i}\right\}=2, \quad\left\{\alpha_{i}, \alpha_{i+1}\right\}=-1, \quad\left\{\alpha_{i}, \alpha_{j}\right\}=0 \quad \text { if }|i-j|>1
$$

Hence

$$
2 \frac{\left\{\alpha_{i}, \alpha_{j}\right\}}{\left\{\alpha_{i}, \alpha_{i}\right\}}=A_{i j} \quad i, j=1, \ldots, l
$$

where $A=\left(A_{i j}\right)$ is the Cartan matrix of type $A_{l}$.

We now consider the action of the fundamental reflections on $V$. We define linear maps $s_{i}: \tilde{V} \rightarrow \tilde{V}$ by

$$
\begin{aligned}
& s_{i}\left(\beta_{i}\right)=\beta_{i+1} \\
& s_{i}\left(\beta_{i+1}\right)=\beta_{i} \\
& s_{i}\left(\beta_{j}\right)=\beta_{j} \quad j \neq i, i+1 .
\end{aligned}
$$

Then we have

$$
s_{i}\left(\alpha_{j}\right)=\alpha_{j}-A_{i j} \alpha_{i} \quad i, j=1, \ldots, l
$$

and so $s_{i}$ restricted to $V$ is the $i$ th fundamental reflection.
We consider the group of transformations of $\tilde{V}$ generated by $s_{1}, \ldots, s_{l}$. Since $s_{i}$ acts on $\tilde{V}$ by permuting $\beta_{i}, \beta_{i+1}$ and fixing the remaining $\beta_{j}$ the group generated by the $s_{i}$ is the group of all permutations of $\beta_{1}, \ldots, \beta_{l+1}$. This group leaves the subspace $V$ invariant and induces on $V$ the Weyl group $W$. Thus we have a surjective homomorphism

$$
S_{l+1} \rightarrow W
$$

whose kernel is trivial. Hence the Weyl group of type $A_{l}$ is isomorphic to the symmetric group $S_{l+1}$.

The full root system $\Phi$ of type $A_{l}$ is the set of vectors of form $w\left(\alpha_{i}\right)$ for all $w \in W$ and all $i$. This is the set

$$
\Phi=\left\{\beta_{i}-\beta_{j} ; i \neq j \quad i, j=1, \ldots, l+1\right\} .
$$

Thus we have $|\Phi|=l(l+1)$ and $\operatorname{dim} L=l+|\Phi|=l(l+2)$.
We shall now show that $\mathfrak{\xi l} l_{l+1}(\mathbb{C})$ is a simple Lie algebra of type $A_{l}$. We discussed this Lie algebra in Section 4.4. In particular we know from Proposition 4.26 that the subalgebra $H$ of diagonal matrices in $L=\mathfrak{\xi l} l_{l+1}(\mathbb{C})$ is a Cartan subalgebra. Moreover by Proposition 4.27

$$
L=H \oplus \sum_{i \neq j} \mathbb{C} E_{i j}
$$

is the Cartan decomposition of $L$ with respect to $H$. By Theorem $4.25 L$ is a simple Lie algebra. By Proposition 4.28 the roots of $L$ are the functions

$$
\left(\begin{array}{cccc}
\lambda_{1} & & & \\
& \cdot & & \\
& & \cdot & \\
& & \cdot & \\
& & & \lambda_{l+1}
\end{array}\right) \rightarrow \lambda_{i}-\lambda_{j} \quad i \neq j
$$

and a system of fundamental roots is given by

$$
\alpha_{i}\left(\begin{array}{cccc}
\lambda_{1} & & & \\
& \cdot & & \\
& & \cdot & \\
& & \cdot & \\
& & & \lambda_{l+1}
\end{array}\right)=\lambda_{i}-\lambda_{i+1}
$$

We can now determine the Cartan matrix $A=\left(A_{i j}\right)$ of $L$. We recall from Proposition 4.22 that the $\alpha_{i}$-chain of roots through $\alpha_{j}$ when $i \neq j$ has the form

$$
\alpha_{j}, \alpha_{i}+\alpha_{j}, \ldots, q \alpha_{i}+\alpha_{j}
$$

where $q=-A_{i j}$. Since we know the roots we can determine the numbers $q$. We have $q=1$ if $i=j-1$ or $j+1$ and $q=0$ otherwise. Thus the Cartan matrix $A$ is the same as the Cartan matrix of type $A_{l}$ in the standard list 6.12. Thus we have proved:

Theorem 8.1 (i) The simple Lie algebra of type $A_{l}$ has dimension $l(l+2)$.
(ii) The Lie algebra $\mathfrak{\xi l}_{l+1}(\mathbb{C})$ of all $(l+1) \times(l+1)$ matrices of trace 0 is simple of type $A_{l}$.

### 8.2 Lie algebras of type $D_{l}$

We recall that the Dynkin diagram of type $D_{l}$ has form


Let $V$ be a real vector space of dimension $l$ and basis $\beta_{1}, \ldots, \beta_{l}$. Let the symmetric scalar product $\{$,$\} be defined by \left\{\beta_{i}, \beta_{j}\right\}=\delta_{i j}$. We define $\alpha_{1}, \ldots, \alpha_{l}$ by

$$
\alpha_{1}=\beta_{1}-\beta_{2}, \quad \alpha_{2}=\beta_{2}-\beta_{3}, \quad \ldots, \quad \alpha_{l-1}=\beta_{l-1}-\beta_{l}, \quad \alpha_{l}=\beta_{l-1}+\beta_{l}
$$

Then we have

$$
\begin{aligned}
& \left\{\alpha_{i}, \alpha_{i}\right\}=2 \quad \text { for all } i \\
& \left\{\alpha_{i}, \alpha_{i+1}\right\}=-1 \quad \text { for } 1 \leq i \leq l-2 \\
& \left\{\alpha_{i}, \alpha_{j}\right\}=0 \quad \text { for } i, j \in\{1, \ldots, l-1\} \text { with }|i-j|>1 \\
& \left\{\alpha_{l-2}, \alpha_{l}\right\}=-1 \\
& \left\{\alpha_{i}, \alpha_{l}\right\}=0 \quad \text { for } i \neq l-2, l
\end{aligned}
$$

It follows that

$$
2 \frac{\left\{\alpha_{i}, \alpha_{j}\right\}}{\left\{\alpha_{i}, \alpha_{i}\right\}}=A_{i j} \quad \text { for all } i, j
$$

and hence that the scalar product $\{$,$\} is a non-zero multiple of the Killing form.$
We now consider the fundamental reflections $s_{i}$ on $V$. For $1 \leq i \leq l-1$ we have

$$
\begin{aligned}
& s_{i}\left(\beta_{i}\right)=\beta_{i+1} \\
& s_{i}\left(\beta_{i+1}\right)=\beta_{i} \\
& s_{i}\left(\beta_{j}\right)=\beta_{j} \quad \text { for } j \neq i, i+1
\end{aligned}
$$

For $i=l$ we have

$$
\begin{aligned}
& s_{l}\left(\beta_{l-1}\right)=-\beta_{l} \\
& s_{l}\left(\beta_{l}\right)=-\beta_{l-1} \\
& s_{l}\left(\beta_{j}\right)=\beta_{j} \quad \text { for } j \neq l-1, l
\end{aligned}
$$

Thus the Weyl group $W$ generated by $s_{1}, \ldots, s_{l}$ has form

$$
w\left(\beta_{i}\right)= \pm \beta_{\sigma(i)} \quad w \in W
$$

for some permutation $\sigma$ of $1, \ldots, l$. Let $w\left(\beta_{i}\right)=\varepsilon_{i} \beta_{\sigma(i)}$. Then an even number of the signs $\varepsilon_{i}$ are equal to -1 . Conversely for any permutation $\sigma$ of $1, \ldots, l$ and any set of signs $\varepsilon_{i}$ with $\prod \varepsilon_{i}=1$ there is an element $w \in W$ acting as above. It follows that the order of the Weyl group of type $D_{l}$ is given by

$$
|W|=2^{l-1} l!
$$

We now consider the root system $\Phi$. The elements of $\Phi$ have form $w\left(\alpha_{i}\right)$ for all $w \in W$ and all $i$. Since $w$ acts on the $\beta_{i}$ by a permutation combined with certain sign changes we obtain

$$
\Phi=\left\{ \pm \beta_{i} \pm \beta_{j} ; i \neq j \in\{1, \ldots, l\}\right\}
$$

All combinations of signs are possible. Hence $|\Phi|=2 l(l-1)$ and so

$$
\operatorname{dim} L=l+|\Phi|=l(2 l-1)
$$

We now wish to describe $L$ as a Lie algebra of matrices. We begin with a lemma which will be useful both for the type being considered and for certain other types also.

Lemma 8.2 Let $M$ be an $n \times n$ matrix over $\mathbb{C}$. Then the set of all $n \times n$ matrices $X$ over $\mathbb{C}$ satisfying

$$
X^{\mathrm{t}} M+M X=O
$$

forms a Lie algebra under Lie multiplication of matrices.
Proof. The set of such matrices $X$ is clearly closed under addition and scalar multiplication. Let $X_{1}, X_{2}$ be matrices satisfying the given condition. Thus we have

$$
X_{1}^{\mathrm{t}} M=-M X_{1}, \quad X_{2}^{\mathrm{t}} M=-M X_{2}
$$

It follows that

$$
\begin{aligned}
{\left[X_{1} X_{2}\right]^{\mathrm{t}} M } & =\left(X_{1} X_{2}-X_{2} X_{1}\right)^{\mathrm{t}} M=X_{2}^{\mathrm{t}} X_{1}^{\mathrm{t}} M-X_{1}^{\mathrm{t}} X_{2}^{\mathrm{t}} M \\
& =-X_{2}^{\mathrm{t}} M X_{1}+X_{1}^{\mathrm{t}} M X_{2}=M X_{2} X_{1}-M X_{1} X_{2} \\
& =-M\left[X_{1} X_{2}\right]
\end{aligned}
$$

Thus the set of such matrices $X$ forms a Lie algebra.
We now consider the special case when $M$ is the $2 l \times 2 l$ matrix

$$
M=\left(\begin{array}{ll}
O & I_{l} \\
I_{l} & O
\end{array}\right)
$$

Then a $2 l \times 2 l$ matrix $X=\left(\begin{array}{ll}X_{11} & X_{12} \\ X_{21} & X_{22}\end{array}\right)$ satisfies $X^{\mathrm{t}} M+M X=O$ if and only if $X_{22}=-X_{11}^{\mathrm{t}}$ and $X_{12}, X_{21}$ are skew-symmetric. Let $L$ be the Lie algebra of all such matrices $X$ and $H$ be the set of diagonal matrices in $L$. The elements of $H$ have form

$$
h=\left(\begin{array}{llllll}
\lambda_{1} & & & & & \\
& \ddots & & & & \\
& & \lambda_{l} & & & \\
& & & -\lambda_{1} & & \\
& & & & \ddots & \\
& & & & & -\lambda_{l}
\end{array}\right)
$$

Let us number the rows and columns $1, \ldots, l,-1, \ldots,-l$. Then we have

$$
L=H \oplus \sum_{\alpha} \mathbb{C} e_{\alpha}
$$

where

$$
e_{\alpha}=\left\{\begin{array}{l}
E_{i j}-E_{-j,-i} \\
-E_{-i,-j}+E_{j i} \quad \text { for } 0<i<j, \\
E_{i,-j}-E_{j,-i} \\
-E_{-i, j}+E_{-j, i}
\end{array}\right.
$$

that is for each pair $i, j$ with $0<i<j$ we have four vectors $e_{\alpha}$ as above. Moreover each of the 1 -dimensional spaces $\mathbb{C} e_{\alpha}$ is a $H$-module, and we have:

$$
\begin{aligned}
& {\left[h, E_{i j}-E_{-j,-i}\right]=\left(\lambda_{i}-\lambda_{j}\right)\left(E_{i j}-E_{-j,-i}\right)} \\
& {\left[h,-E_{-i,-j}+E_{j i}\right]=\left(\lambda_{j}-\lambda_{i}\right)\left(-E_{-i,-j}+E_{j i}\right)} \\
& {\left[h, E_{i,-j}-E_{j,-i}\right]=\left(\lambda_{i}-\lambda_{j}\right)\left(E_{i,-j}-E_{j,-i}\right)} \\
& {\left[h,-E_{-i, j}+E_{-j, i}\right]=\left(-\lambda_{i}-\lambda_{j}\right)\left(-E_{-i, j}+E_{-j, i}\right) .}
\end{aligned}
$$

We write $\left[h e_{\alpha}\right]=\alpha(h) e_{\alpha}$ for all such $e_{\alpha}$.
Now the argument of Proposition 7.34 shows that $H$ is a Cartan subalgebra of $L$. The decomposition

$$
L=H \oplus \sum_{\alpha} \mathbb{C} e_{\alpha}
$$

is then the Cartan decomposition of $L$ with respect to $H$.
We next verify that $L$ is semisimple. Suppose not. Then $L$ has a non-zero abelian ideal $I$. Since $[H I] \subset I$ we may regard $I$ as a $H$-module and consider the decomposition of $I$ into weight spaces with respect to $H$. This gives

$$
I=(H \cap I) \oplus \sum_{\alpha}\left(\mathbb{C} e_{\alpha} \cap I\right)
$$

just as in the proof of Proposition 7.33. Suppose if possible that $\mathbb{C} e_{\alpha} \cap I \neq O$ for some $\alpha$. Then we have $e_{\alpha} \in I$. We then define $h_{\alpha}$ by $h_{\alpha}=\left[e_{\alpha} e_{-\alpha}\right]$ and observe that $\left[h_{\alpha} e_{\alpha}\right]=2 e_{\alpha}$. Then $e_{\alpha}, h_{\alpha} \in I$ and we have a contradiction to the fact that $I$ is abelian. Hence $\mathbb{C} e_{\alpha} \cap I=O$ for all $\alpha$ and so $I \subset H$. Let $x \in I$. Then $\left[x e_{\alpha}\right]=\alpha(x) e_{\alpha} \in I$ so $\alpha(x)=0$. This holds for all $\alpha$ and so $x=0$. Thus $I=O$, which gives a contradiction. Hence $L$ is semisimple.

We now know that the functions $\alpha: H \rightarrow \mathbb{C}$ given above are the roots of $L$ with respect to $H$. A system of fundamental roots is given by

$$
\begin{aligned}
& \alpha_{1}(h)=\lambda_{1}-\lambda_{2} \\
& \alpha_{2}(h)=\lambda_{2}-\lambda_{3} \\
& \vdots \\
& \alpha_{l-1}(h)=\lambda_{l-1}-\lambda_{l} \\
& \alpha_{l}(h)=\lambda_{l-1}+\lambda_{l}
\end{aligned}
$$

since all the other roots are integral combinations of these with coefficients all non-negative or all non-positive.

We now determine the Cartan matrix of $L$. Let the $\alpha_{i}$-chain of roots through $\alpha_{j}$ for $i \neq j$ be

$$
\alpha_{j}, \quad \alpha_{i}+\alpha_{j}, \quad \ldots, \quad q \alpha_{i}+\alpha_{j}
$$

Then $A_{i j}=-q$ by Proposition 4.22. Since we know the roots we can find the number $q$ and hence $A_{i j}$ for each $i \neq j$. This gives us the Cartan matrix $A=\left(A_{i j}\right)$ of type $D_{l}$ on the standard list 6.12.

Finally we note that since this Cartan matrix is indecomposable the Lie algebra $L$ must be simple by Corollary 6.15 . Thus we have proved the following result:

Theorem 8.3 (i) The simple Lie algebra of type $D_{l}$ has dimension $l(2 l-1)$.
(ii) The Lie algebra of all $2 l \times 2 l$ matrices $X$ satisfying $X^{t} M+M X=0$ where

$$
M=\left(\begin{array}{ll}
O & I_{l} \\
I_{l} & O
\end{array}\right)
$$

is simple of type $D_{l}$ when $l \geq 4$.

### 8.3 Lie algebras of type $B_{l}$

We recall that the Dynkin diagram of type $B_{l}$ has form


Let $V$ be a real vector space of dimension $l$ with basis $\beta_{1}, \ldots, \beta_{l}$. Let the scalar product $\{$,$\} on V$ be defined by $\left\{\beta_{i}, \beta_{j}\right\}=\delta_{i j}$. We define $\alpha_{1}, \ldots, \alpha_{l} \in V$ by

$$
\alpha_{1}=\beta_{1}-\beta_{2}, \quad \alpha_{2}=\beta_{2}-\beta_{3}, \quad \ldots, \quad \alpha_{l-1}=\beta_{l-1}-\beta_{l}, \quad \alpha_{l}=\beta_{l}
$$

Then we have

$$
\begin{array}{lc}
\left\{\alpha_{i}, \alpha_{i}\right\}=2 & \text { for } 1 \leq i \leq l-1 \\
\left\{\alpha_{l}, \alpha_{l}\right\}=1 & \\
\left\{\alpha_{i}, \alpha_{i+1}\right\}=-1 & \text { for } 1 \leq i \leq l-1 \\
\left\{\alpha_{i}, \alpha_{j}\right\}=0 & \text { if }|i-j|>1
\end{array}
$$

It follows that

$$
2 \frac{\left\{\alpha_{i}, \alpha_{j}\right\}}{\left\{\alpha_{i}, \alpha_{i}\right\}}=A_{i j} \quad \text { for all } i, j
$$

where $A=\left(A_{i j}\right)$ is the Cartan matrix of type $B_{l}$ on the standard list 6.12. Thus the scalar product $\{$,$\} is a non-zero multiple of the Killing form.$

We now consider the fundamental reflections $s_{i}$ on $V$. We have, for $1 \leq$ $i \leq l-1$,

$$
\begin{aligned}
& s_{i}\left(\beta_{i}\right)=\beta_{i+1} \\
& s_{i}\left(\beta_{i+1}\right)=\beta_{i} \\
& s_{i}\left(\beta_{j}\right)=\beta_{j} \quad j \neq i, i+1 .
\end{aligned}
$$

For $i=l$ we have

$$
\begin{aligned}
& s_{l}\left(\beta_{l}\right)=-\beta_{l} \\
& s_{l}\left(\beta_{i}\right)=\beta_{i} \quad i \neq l .
\end{aligned}
$$

Thus the Weyl group $W$ generated by $s_{1}, \ldots, s_{l}$ consists of elements $w$ of the form

$$
w\left(\beta_{i}\right)= \pm \beta_{\sigma(i)}
$$

for some permutation $\sigma$ of $1, \ldots, l$. Let $w\left(\beta_{i}\right)=\varepsilon_{i} \beta_{\sigma(i)}$. Then, given any permutation $\sigma$ of $1, \ldots, l$ and any set of signs $\varepsilon_{i} \in\{1,-1\}$ there is an element $w \in W$ such that $w\left(\beta_{i}\right)=\varepsilon_{i} \beta_{\sigma(i)}$ for all $i$. Thus the order of the Weyl group $W$ of type $B_{l}$ is

$$
|W|=2^{l} l!.
$$

We now consider the root system $\Phi$. The elements of $\Phi$ have form $w\left(\alpha_{i}\right)$ for all $w \in W$ and all $i$. Since $w$ acts on the $\beta_{i}$ by means of a permutation combined with sign changes we obtain

$$
\Phi=\left\{ \pm \beta_{i} \pm \beta_{j} \quad i \neq j ; \pm \beta_{i}\right\}
$$

All combinations of signs are possible. Thus we have $|\Phi|=2 l^{2}$ and so

$$
\operatorname{dim} L=l+|\Phi|=l(2 l+1)
$$

We shall now describe $L$ as a Lie algebra of matrices. We use Lemma 8.2 and this time we take the $(2 l+1) \times(2 l+1)$ matrix $M$ given by

$$
M=\left(\begin{array}{cccc}
2 & 0 & \ldots & 0 \\
0 & O & & I_{l} \\
\vdots & & & \\
0 & I_{l} & & O
\end{array}\right) .
$$

Let $L$ be the Lie algebra of all $(2 l+1) \times(2 l+1)$ matrices $X$ satisfying the condition

$$
X^{\mathrm{t}} M+M X=O
$$

We consider $X$ as a block matrix

$$
\begin{array}{ccc}
\left(\begin{array}{ccc}
X_{00} & X_{01} & X_{02} \\
X_{10} & X_{11} & X_{12} \\
X_{20} & X_{21} & X_{22}
\end{array}\right) & l \\
l & l & l
\end{array}
$$

Then $X$ satisfies $X^{\mathrm{t}} M+M X=O$ if and only if $X_{22}=-X_{11}^{\mathrm{t}}, X_{12}$ and $X_{21}$ are skew-symmetric, $X_{10}=-2 X_{02}^{\mathrm{t}}, X_{20}=-2 X_{01}^{\mathrm{t}}$ and $X_{00}=0$.

Let $H$ be the set of diagonal matrices in $L$. The elements of $H$ have form

$$
h=\left(\begin{array}{lllllll}
0 & & & & & & \\
& \lambda_{1} & & & & & \\
& & \ddots & & & & \\
& & & \lambda_{1} & & & \\
& & & & -\lambda_{1} & & \\
& & & & & \ddots & \\
& & & & & & -\lambda_{l}
\end{array}\right) .
$$

We number the rows and columns $0,1, \ldots, l,-1, \ldots,-l$. Then we have

$$
L=H \oplus \sum_{\alpha} \mathbb{C} e_{\alpha}
$$

where

$$
e_{\alpha}= \begin{cases}E_{i j}-E_{-j,-i} & \\ -E_{-i,-j}+E_{j i} & \\ E_{i,-j}-E_{j,-i} & \text { for } 0<i<j \\ -E_{-i, j}+E_{-j, i} & \\ 2 E_{i 0}-E_{0,-i} & \text { for } 0<i \\ -2 E_{-i 0}+E_{0 i} & \end{cases}
$$

Each of the 1-dimensional spaces $\mathbb{C}_{e_{\alpha}}$ is an $H$-module and we have

$$
\begin{aligned}
& {\left[h, E_{i j}-E_{-j,-i}\right]=\left(\lambda_{i}-\lambda_{j}\right)\left(E_{i j}-E_{-j,-i}\right)} \\
& {\left[h,-E_{-i,-j}+E_{j i}\right]=\left(\lambda_{j}-\lambda_{i}\right)\left(-E_{-i,-j}+E_{j i}\right)} \\
& {\left[h, E_{i,-j}-E_{j,-i}\right]=\left(\lambda_{i}+\lambda_{j}\right)\left(E_{i,-j}-E_{j,-i}\right)} \\
& {\left[h,-E_{-i, j}+E_{-j, i}\right]=\left(-\lambda_{i}-\lambda_{j}\right)\left(-E_{-i, j}+E_{-j, i}\right)} \\
& {\left[h, 2 E_{i 0}-E_{0,-i}\right]=\lambda_{i}\left(2 E_{i 0}-E_{0,-i}\right)} \\
& {\left[h,-2 E_{-i, 0}+E_{0 i}\right]=-\lambda_{i}\left(-2 E_{-i, 0}+E_{0 i}\right) .}
\end{aligned}
$$

We write $\left[h e_{\alpha}\right]=\alpha(h) e_{\alpha}$ for all such $\alpha$.
We now show that $H$ is a Cartan subalgebra of $L ; L=H \oplus \sum_{\alpha} \mathbb{C} e_{\alpha}$ is the Cartan decomposition of $L$ with respect to $H$, and $L$ is semisimple. These facts can be proved in exactly the same way as that used in Section 8.2 for type $D_{l}$.

We now know that the functions $\alpha: H \rightarrow \mathbb{C}$ given above are the roots of $L$ with respect to $H$. A system of fundamental roots is given by

$$
\begin{aligned}
& \alpha_{1}(h)=\lambda_{1}-\lambda_{2} \\
& \alpha_{2}(h)=\lambda_{2}-\lambda_{3} \\
& \quad \vdots \\
& \alpha_{l-1}(h)=\lambda_{l-1}-\lambda_{l} \\
& \alpha_{l}(h)=\lambda_{l}
\end{aligned}
$$

since all the other roots are integral combinations of these with coefficients all non-negative or all non-positive.

We can now determine the Cartan integers $A_{i j}$. Let $\alpha_{j}, \alpha_{i}+\alpha_{j}, \ldots, q \alpha_{i}+\alpha_{j}$ be the $\alpha_{i}$-chain of roots through $\alpha_{j}$ for $i \neq j$. By Proposition 4.22 we have $A_{i j}=-q$ and so the Cartan matrix $A=\left(A_{i j}\right)$ can be determined. This turns
out to be the Cartan matrix of type $B_{l}$ on the standard list 6.12 . Finally we observe that $L$ must be a simple Lie algebra by Corollary 6.15, since its Cartan matrix is indecomposable. Thus we have

Theorem 8.4 (i) The simple Lie algebra of type $B_{l}$ has dimension $l(2 l+1)$.
(ii) The Lie algebra of all $(2 l+1) \times(2 l+1)$ matrices $X$ satisfying $X^{\mathrm{t}} M+$ $M X=O$ where

$$
M=\left(\begin{array}{cccc}
2 & 0 & \cdots & 0 \\
0 & O & & I_{l} \\
\vdots & & & \\
0 & I_{l} & & O
\end{array}\right)
$$

is simple of type $B_{l}$ when $l \geq 2$.

### 8.4 Lie algebras of type $C_{l}$

We recall that the Dynkin diagram of type $C_{l}$ has form


Let $V$ be a real vector space of dimension $l$ with basis $\beta_{1}, \ldots, \beta_{l}$. Let the scalar product $\{$,$\} on V$ be defined by $\left\{\beta_{i}, \beta_{j}\right\}=\delta_{i j}$. We define $\alpha_{1}, \ldots, \alpha_{l} \in V$ by

$$
\alpha_{1}=\beta_{1}-\beta_{2}, \quad \alpha_{2}=\beta_{2}-\beta_{3}, \quad \ldots, \quad \alpha_{l-1}=\beta_{l-1}-\beta_{l}, \quad \alpha_{l}=2 \beta_{l}
$$

Then we have

$$
\begin{aligned}
& \left\{\alpha_{i}, \alpha_{i}\right\}=2 \quad \text { for } 1 \leq i \leq l-1 \\
& \left\{\alpha_{l}, \alpha_{l}\right\}=4 \\
& \left\{\alpha_{i}, \alpha_{i+1}\right\}=-1 \quad \text { for } 1 \leq i \leq l-2 \\
& \left\{\alpha_{l-1}, \alpha_{l}\right\}=-2 \\
& \left\{\alpha_{i}, \alpha_{j}\right\}=0 \quad \text { for }|i-j|>1
\end{aligned}
$$

It follows that

$$
2 \frac{\left\{\alpha_{i}, \alpha_{j}\right\}}{\left\{\alpha_{i}, \alpha_{i}\right\}}=A_{i j} \quad \text { for all } i, j
$$

where $A=\left(A_{i j}\right)$ is the Cartan matrix of type $C_{l}$ on the standard list 6.12. Thus the scalar product $\{$,$\} is a non-zero multiple of the Killing form.$

The fundamental reflections $s_{1}, \ldots, s_{l}$ act on $\beta_{1}, \ldots, \beta_{l}$ in exactly the same manner as in type $B_{l}$, considered in Section 8.3. Thus we have

$$
|W|=2^{l} l!
$$

as in Section 8.3 and each $w \in W$ acts on the $\beta_{i}$ by means of a permutation combined with sign changes. Both the permutation and the sign changes can be chosen arbitrarily. Thus we obtain the root system $\Phi$ as the set of all vectors of form $w\left(\alpha_{i}\right)$ for all $w \in W$ and all $i$. Thus

$$
\Phi=\left\{ \pm \beta_{i} \pm \beta_{j} \quad i \neq j ; \pm 2 \beta_{i}\right\}
$$

All combinations of signs are possible. Thus we have $|\Phi|=2 l^{2}$ and

$$
\operatorname{dim} L=l+|\Phi|=l(2 l+1)
$$

We next describe $L$ as a Lie algebra of matrices. Again we use Lemma 8.2. This time we take the $2 l \times 2 l$ matrix $M$ given by

$$
M=\left(\begin{array}{cc}
O & I_{l} \\
-I_{l} & O
\end{array}\right)
$$

Let $L$ be the Lie algebra of all $2 l \times 2 l$ matrices satisfying the condition

$$
X^{\mathrm{t}} M+M X=O
$$

Let $X=\left(\begin{array}{ll}X_{11} & X_{12} \\ X_{21} & X_{22}\end{array}\right)$. Then $X$ lies in $L$ if and only if $X_{22}=-X_{11}^{\mathrm{t}}$ and $X_{12}, X_{21}$ are symmetric.

Let $H$ be the set of diagonal matrices in $L$. The elements of $H$ have form

$$
h=\left(\begin{array}{llllll}
\lambda_{1} & & & & & \\
& \ddots & & & & \\
& & \lambda_{l} & & & \\
& & & -\lambda_{1} & & \\
& & & & \ddots & \\
& & & & & -\lambda_{l}
\end{array}\right)
$$

We number the rows and columns $1, \ldots, l,-1, \ldots,-l$. Then we have

$$
L=H \oplus \sum_{\alpha} \mathbb{C} e_{\alpha}
$$

where

$$
e_{\alpha}=\left\{\begin{array}{l}
E_{i j}-E_{-j,-i} \\
-E_{-i,-j}+E_{j i} \quad 0<i<j \\
E_{i,-j}+E_{j,-i} \\
E_{-i, j}+E_{-j, i} \\
E_{i,-i} \quad 0<i . \\
E_{-i, i}
\end{array}\right.
$$

Each of the 1 -dimensional spaces $\mathbb{C} e_{\alpha}$ is a $H$-module. We have

$$
\begin{aligned}
& {\left[h, E_{i j}-E_{-j,-i}\right]=\left(\lambda_{i}-\lambda_{j}\right)\left(E_{i j}-E_{-j,-i}\right)} \\
& {\left[h,-E_{-i,-j}+E_{j i}\right]=\left(\lambda_{j}-\lambda_{i}\right)\left(-E_{-i,-j}+E_{j i}\right)} \\
& {\left[h, E_{i,-j}+E_{j,-i}\right]=\left(\lambda_{i}+\lambda_{j}\right)\left(E_{i,-j}+E_{j,-i}\right)} \\
& {\left[h, E_{-i, j}+E_{-j, i}\right]=\left(-\lambda_{i}-\lambda_{j}\right)\left(E_{-i, j}+E_{-j, i}\right)} \\
& {\left[h, E_{i,-i}\right]=2 \lambda_{i} E_{i,-i}} \\
& {\left[h, E_{-i, i}\right]=-2 \lambda_{i} E_{-i, i} .}
\end{aligned}
$$

We write $\left[h e_{\alpha}\right]=\alpha(h) e_{\alpha}$ for all such $\alpha$.
We observe that $H$ is a Cartan subalgebra of $L$, that $L=H \oplus \sum_{\alpha} \mathbb{C} e_{\alpha}$ is the Cartan decomposition of $L$ with respect to $H$, and that the Lie algebra $L$ is semisimple, using the same arguments as given in Section 8.2 for type $D_{l}$.

The functions $\alpha: H \rightarrow \mathbb{C}$ given above are the roots of $L$ with respect to $H$. A system of fundamental roots is given by

$$
\begin{aligned}
& \alpha_{1}(h)=\lambda_{1}-\lambda_{2} \\
& \alpha_{2}(h)=\lambda_{2}-\lambda_{3} \\
& \vdots \\
& \alpha_{l-1}(h)=\lambda_{l-1}-\lambda_{l} \\
& \alpha_{l}(h)=2 \lambda_{l}
\end{aligned}
$$

since all the other roots are integral combinations of these with coefficients all non-negative or all non-positive.

We can now determine the Cartan integers $A_{i j}$. Let $\alpha_{j}, \alpha_{i}+\alpha_{j}, \ldots, q \alpha_{i}+\alpha_{j}$ be the $\alpha_{i}$-chain of roots through $\alpha_{j}$, for $i \neq j$. Then $A_{i j}=-q$. The Cartan matrix $A=\left(A_{i j}\right)$ determined in this way turns out to be the Cartan matrix of
type $C_{l}$ on the standard list 6.12 . Finally $L$ is a simple Lie algebra, since its Cartan matrix is indecomposable. Thus we have

Theorem 8.5 (i) The simple Lie algebra of type $C_{l}$ has dimension $l(2 l+1)$.
(ii) The Lie algebra of all $2 l \times 2 l$ matrices $X$ satisfying $X^{\mathrm{t}} M+M X=O$ where

$$
M=\left(\begin{array}{cc}
O & I_{l} \\
-I_{l} & O
\end{array}\right)
$$

is simple of type $C_{l}$ when $l \geq 3$.
The Lie algebras of type $A_{l}, B_{l}, C_{l}$ or $D_{l}$ are called the simple Lie algebras of classical type. The remaining simple Lie algebras $E_{6}, E_{7}, E_{8}, F_{4}$, $G_{2}$ are called the exceptional simple Lie algebras. We now determine the dimensions and root systems of the exceptional Lie algebras.

### 8.5 Lie algebras of type $G_{2}$

The Dynkin diagram of type $G_{2}$ is

and the corresponding Cartan matrix is

$$
\left(\begin{array}{cc}
2 & -1 \\
-3 & 2
\end{array}\right)
$$

Let $\alpha_{1}, \alpha_{2}$ be the fundamental roots in a root system of type $G_{2}$. Then we have

$$
\begin{array}{ll}
s_{1}\left(\alpha_{1}\right)=-\alpha_{1} & s_{2}\left(\alpha_{1}\right)=\alpha_{1}+3 \alpha_{2} \\
s_{1}\left(\alpha_{2}\right)=\alpha_{1}+\alpha_{2} & s_{2}\left(\alpha_{2}\right)=-\alpha_{2}
\end{array}
$$

and $W=\left\langle s_{1}, s_{2}\right\rangle$. Thus each root in $\Phi$ is obtained from $\alpha_{1}$ or $\alpha_{2}$ by applying $s_{1}, s_{2}$ alternately. Now we have

$$
\begin{aligned}
& \alpha_{1} \rightarrow-\alpha_{1} \rightarrow \underset{s_{1}}{\rightarrow}-\alpha_{1}-3 \alpha_{2} \rightarrow-2 \alpha_{s_{1}}-3 \alpha_{2} \\
& \alpha_{1} \rightarrow \alpha_{s_{2}}+3 \alpha_{2} \rightarrow 2 \alpha_{1}+3 \alpha_{s_{1}} \\
& \alpha_{2} \rightarrow \alpha_{s_{1}}+\alpha_{2} \rightarrow \alpha_{s_{2}}+2 \alpha_{2} \\
& \alpha_{2} \rightarrow-\alpha_{s_{2}} \rightarrow \alpha_{s_{1}}-\alpha_{1}-\alpha_{2} \rightarrow-\alpha_{s_{2}}-2 \alpha_{2}
\end{aligned}
$$

and

$$
\begin{aligned}
& s_{2}\left(-2 \alpha_{1}-3 \alpha_{2}\right)=-2 \alpha_{1}-3 \alpha_{2} \\
& s_{2}\left(2 \alpha_{1}+3 \alpha_{2}\right)=2 \alpha_{1}+3 \alpha_{2} \\
& s_{1}\left(\alpha_{1}+2 \alpha_{2}\right)=\alpha_{1}+2 \alpha_{2} \\
& s_{1}\left(-\alpha_{1}-2 \alpha_{2}\right)=-\alpha_{1}-2 \alpha_{2}
\end{aligned}
$$

Thus all the vectors in the above sequences are roots, and we do not obtain new vectors by continuing the sequences further. Hence

$$
\begin{gathered}
\Phi=\left\{\begin{array}{lllll}
\alpha_{1}, & \alpha_{2}, & \alpha_{1}+\alpha_{2}, \quad \alpha_{1}+2 \alpha_{2}, & \alpha_{1}+3 \alpha_{2}, & 2 \alpha_{1}+3 \alpha_{2},
\end{array} \quad-\alpha_{1},\right. \\
-\alpha_{2},-\alpha_{1}-\alpha_{2}, \quad-\alpha_{1}-2 \alpha_{2}, \quad-\alpha_{1}-3 \alpha_{2}, \\
\left.-2 \alpha_{1}-3 \alpha_{2}\right\}
\end{gathered}
$$

Thus we have $|\Phi|=12$ and $\operatorname{dim} L=14$. Hence we have proved

Theorem 8.6 The simple Lie algebra of type $G_{2}$ has dimension 14 .
Figures $8.1,8.2$ and 8.3 compare the simple root systems of types $A_{2}, B_{2}$ and $G_{2}$.


Figure 8.1 Simple root system of type $A_{2}$


Figure 8.2 Simple root system of type $B_{2}$


Figure 8.3 Simple root system of type $G_{2}$

### 8.6 Lie algebras of type $F_{4}$

The Dynkin diagram of type $F_{4}$ is

and the corresponding Cartan matrix is

$$
\left(\begin{array}{cccc}
2 & -1 & 0 & 0 \\
-1 & 2 & -1 & 0 \\
0 & -2 & 2 & -1 \\
0 & 0 & -1 & 2
\end{array}\right)
$$

Let $V$ be a real vector space with $\operatorname{dim} V=4$ and $\beta_{1}, \beta_{2}, \beta_{3}, \beta_{4}$ be a basis of $V$. Let the scalar product $\{$,$\} on V$ be defined by $\left\{\beta_{i}, \beta_{j}\right\}=\delta_{i j}$. We define $\alpha_{1}, \alpha_{2}, \alpha_{3}, \alpha_{4} \in V$ by
$\alpha_{1}=\beta_{1}-\beta_{2} \quad \alpha_{2}=\beta_{2}-\beta_{3} \quad \alpha_{3}=\beta_{3} \quad \alpha_{4}=\frac{1}{2}\left(-\beta_{1}-\beta_{2}-\beta_{3}+\beta_{4}\right)$.
Then we have

$$
\begin{aligned}
& \left\{\alpha_{1}, \alpha_{1}\right\}=\left\{\alpha_{2}, \alpha_{2}\right\}=2 \\
& \left\{\alpha_{3}, \alpha_{3}\right\}=\left\{\alpha_{4}, \alpha_{4}\right\}=1 \\
& \left\{\alpha_{1}, \alpha_{2}\right\}=\left\{\alpha_{2}, \alpha_{3}\right\}=-1 \\
& \left\{\alpha_{3}, \alpha_{4}\right\}=-\frac{1}{2} \\
& \left\{\alpha_{i}, \alpha_{j}\right\}=0 \quad \text { if }|i-j|>1
\end{aligned}
$$

It follows that

$$
2 \frac{\left\{\alpha_{i}, \alpha_{j}\right\}}{\left\{\alpha_{i}, \alpha_{i}\right\}}=A_{i j} \quad \text { for all } i, j
$$

Thus the scalar product $\{$,$\} is a non-zero multiple of the Killing form. We$ consider the action of the corresponding fundamental reflections $s_{1}, s_{2}, s_{3}, s_{4}$. We have

$$
\begin{array}{llll}
s_{1}\left(\beta_{1}\right)=\beta_{2}, & s_{1}\left(\beta_{2}\right)=\beta_{1}, & s_{1}\left(\beta_{3}\right)=\beta_{3}, & s_{1}\left(\beta_{4}\right)=\beta_{4} \\
s_{2}\left(\beta_{1}\right)=\beta_{1}, & s_{2}\left(\beta_{2}\right)=\beta_{3}, & s_{2}\left(\beta_{3}\right)=\beta_{2}, & s_{2}\left(\beta_{4}\right)=\beta_{4} \\
s_{3}\left(\beta_{1}\right)=\beta_{1}, & s_{3}\left(\beta_{2}\right)=\beta_{2}, & s_{3}\left(\beta_{3}\right)=-\beta_{3}, & s_{3}\left(\beta_{4}\right)=\beta_{4}
\end{array}
$$

We consider the subgroup $\left\langle s_{1}, s_{2}, s_{3}\right\rangle$ of the Weyl group $W$ generated by $s_{1}$, $s_{2}, s_{3}$. Elements in this subgroup all fix $\beta_{4}$ but act on $\beta_{1}, \beta_{2}, \beta_{3}$ by means of a permutation combined with sign changes. Thus $w\left(\beta_{i}\right)=\varepsilon_{i} \beta_{\sigma(i)}$ for $i=1,2,3$.

Moreover each permutation $\sigma$ and each choice of signs $\varepsilon_{i}$ arise in this way. Applying the elements of this subgroup of $W$ to $\alpha_{1}, \alpha_{2}, \alpha_{3}, \alpha_{4}$ we see that the vectors

$$
\begin{aligned}
& \pm \beta_{i} \quad 1 \leq i \leq 3 \\
& \pm \beta_{i} \pm \beta_{j} \quad i \neq j \quad 1 \leq i, j \leq 3 \\
& \frac{1}{2}\left( \pm \beta_{1} \pm \beta_{2} \pm \beta_{3} \pm \beta_{4}\right)
\end{aligned}
$$

all lie in $\Phi$. We next consider the action of $s_{4}$. We have

$$
\begin{aligned}
& s_{4}\left(\beta_{1}\right)=\frac{1}{2}\left(\beta_{1}-\beta_{2}-\beta_{3}+\beta_{4}\right) \\
& s_{4}\left(\beta_{2}\right)=\frac{1}{2}\left(-\beta_{1}+\beta_{2}-\beta_{3}+\beta_{4}\right) \\
& s_{4}\left(\beta_{3}\right)=\frac{1}{2}\left(-\beta_{1}-\beta_{2}+\beta_{3}+\beta_{4}\right) \\
& s_{4}\left(\beta_{4}\right)=\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}\right) .
\end{aligned}
$$

Since $s_{4}^{2}=1$ we have $s_{4}\left(\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}\right)\right)=\beta_{4}$. Hence $\beta_{4} \in \Phi$. We also have $s_{4}\left(\beta_{1}+\beta_{2}\right)=-\beta_{3}+\beta_{4}$. Hence $-\beta_{3}+\beta_{4} \in \Phi$. Thus, applying further elements of the subgroup $\left\langle s_{1}, s_{2}, s_{3}\right\rangle$ we see that the vectors

$$
\begin{aligned}
& \pm \beta_{i} \quad 1 \leq i \leq 4 \\
& \pm \beta_{i} \pm \beta_{j} \quad i \neq j \quad 1 \leq i, j \leq 4 \\
& \frac{1}{2}\left( \pm \beta_{1} \pm \beta_{2} \pm \beta_{3} \pm \beta_{4}\right)
\end{aligned}
$$

all lie in $\Phi$, where the choice of signs is arbitrary.
We show this set of vectors is the whole of $\Phi$. To do so it is sufficient to show that the set is invariant under $s_{1}, s_{2}, s_{3}, s_{4}$. The set is clearly invariant under $s_{1}, s_{2}, s_{3}$ because of the simple action of these reflections on $\beta_{1}, \beta_{2}, \beta_{3}$, $\beta_{4}$ described above. Thus it is sufficient to show the set is invariant under $s_{4}$. Now the action of $s_{4}$ given above shows that

$$
s_{4}\left( \pm \beta_{i}\right)=\frac{1}{2}\left(\varepsilon_{1} \beta_{1}+\varepsilon_{2} \beta_{2}+\varepsilon_{3} \beta_{3}+\varepsilon_{4} \beta_{4}\right)
$$

where $\varepsilon_{i} \in\{1,-1\}$ and $\prod \varepsilon_{i}=1$. Thus $s_{4}$ transforms vectors $\pm \beta_{i}$ into the given set, giving as images vectors $\frac{1}{2} \sum \varepsilon_{i} \beta_{i}$ with $\Pi \varepsilon_{i}=1$. Since there are eight such vectors they all appear as vectors $s_{4}\left( \pm \beta_{i}\right)$. Since $s_{4}^{2}=1$ we deduce that all vectors $\frac{1}{2} \sum \varepsilon_{i} \beta_{i}$ with $\prod \varepsilon_{i}=1$ are transformed by $s_{4}$ into the given set.

The formulae for $s_{4}\left(\beta_{i}\right)$ also show that, for all $i \neq j, s_{4}\left( \pm \beta_{i} \pm \beta_{j}\right)$ has form $\pm \beta_{k} \pm \beta_{l}$ for certain $k \neq l$. Thus $s_{4}$ transforms vectors $\pm \beta_{i} \pm \beta_{j}, i \neq j$, into the given set.
It remains to show that $s_{4}$ transforms all vectors $\frac{1}{2} \sum \varepsilon_{i} \beta_{i}$ with $\prod \varepsilon_{i}=$ -1 into the given set. We may clearly assume $\varepsilon_{4}=1$. There are four such vectors. One of them is $\alpha_{4}$ and we have $s_{4}\left(\alpha_{4}\right)=-\alpha_{4}$. The other three are all orthogonal to $\alpha_{4}$ and so are transformed into themselves by $s_{4}$.

Thus the given set of vectors is invariant under $s_{1}, s_{2}, s_{3}, s_{4}$ so is the whole of $\Phi$. Thus we have

$$
\begin{aligned}
\Phi= & \left\{ \pm \beta_{i} \quad 1 \leq i \leq 4\right. \\
& \pm \beta_{i} \pm \beta_{j} \quad i \neq j \quad 1 \leq i, j \leq 4 \\
& \left.\frac{1}{2}\left( \pm \beta_{1} \pm \beta_{2} \pm \beta_{3} \pm \beta_{4}\right)\right\} .
\end{aligned}
$$

In particular we have $|\Phi|=48$, hence $\operatorname{dim} L=52$. Thus we have

Theorem 8.7 The simple Lie algebra of type $F_{4}$ has dimension 52.

We observe that the roots of $F_{4}$ are of two different lengths. There are 24 short roots and 24 long roots. The short roots are $\pm \beta_{i}$ and $\frac{1}{2}\left( \pm \beta_{1} \pm \beta_{2} \pm \beta_{3} \pm \beta_{4}\right)$. The long roots are $\pm \beta_{i} \pm \beta_{j}$.

### 8.7 Lie algebras of types $E_{6}, E_{7}, E_{8}$

We now consider the simple Lie algebra of type $E_{8}$. Its Dynkin diagram is


Let $V$ be a real vector space with $\operatorname{dim} V=8$ and with basis $\beta_{i} i=1, \ldots, 8$. Let the scalar product $\{$,$\} on V$ be defined by $\left\{\beta_{i}, \beta_{j}\right\}=\delta_{i j}$. We wish to find a fundamental system of roots of type $E_{8}$ in $V$. We note that if the vertex 8 is removed from the Dynkin diagram we obtain a Dynkin diagram of type $D_{7}$. This indicates how the first seven vectors in the fundamental system should be chosen. The last one is chosen to be linearly independent of the others and
to satisfy the appropriate conditions relating to the scalar product. Thus we define $\alpha_{1}, \ldots, \alpha_{8} \in V$ by:

$$
\begin{aligned}
\alpha_{i} & =\beta_{i}-\beta_{i+1} \quad 1 \leq i \leq 6 \\
\alpha_{7} & =\beta_{6}+\beta_{7} \\
\alpha_{8} & =-\frac{1}{2} \sum_{i=1}^{8} \beta_{i} .
\end{aligned}
$$

Then we have

$$
\begin{aligned}
& \left\{\alpha_{i}, \alpha_{i}\right\}=2 \quad \text { for } 1 \leq i \leq 8 \\
& \left\{\alpha_{i}, \alpha_{i+1}\right\}=-1 \quad \text { for } 1 \leq i \leq 5 \\
& \left\{\alpha_{5}, \alpha_{7}\right\}=-1 \\
& \left\{\alpha_{7}, \alpha_{8}\right\}=-1 \\
& \left\{\alpha_{i}, \alpha_{j}\right\}=0 \quad \text { for all other pairs } i, j
\end{aligned}
$$

It follows that

$$
2 \frac{\left\{\alpha_{i}, \alpha_{j}\right\}}{\left\{\alpha_{i}, \alpha_{i}\right\}}=A_{i j}
$$

where $A=\left(A_{i j}\right)$ is the Cartan matrix of type $E_{8}$ on the standard list.
In order to obtain the remaining roots we consider the action of the fundamental reflections $s_{1}, \ldots, s_{8}$. We have

$$
\begin{aligned}
& s_{i}\left(\beta_{i}\right)=\beta_{i+1} \\
& s_{i}\left(\beta_{i+1}\right)=\beta_{i} \\
& s_{i}\left(\beta_{j}\right)=\beta_{j} \quad \text { for } j \neq i, i+1
\end{aligned}
$$

when $1 \leq i \leq 6$. Thus the subgroup of the Weyl group $W$ generated by $s_{1}, \ldots, s_{6}$ will give all permutations of $\beta_{1}, \ldots, \beta_{7}$ and will fix $\beta_{8}$. The fundamental reflection $s_{7}$ acts by:

$$
\begin{aligned}
& s_{7}\left(\beta_{6}\right)=-\beta_{7} \\
& s_{7}\left(\beta_{7}\right)=-\beta_{6} \\
& s_{7}\left(\beta_{i}\right)=\beta_{i} \quad i \neq 6,7
\end{aligned}
$$

Thus the subgroup of $W$ generated by $s_{1}, \ldots, s_{7}$ will act on $\beta_{1}, \ldots, \beta_{7}$ by permutations and sign changes, and will fix $\beta_{8}$. Moreover the number of sign changes will be even, and any permutation of $\beta_{1}, \ldots, \beta_{7}$ combined with any even number of sign changes will arise in this way.

It is then clear that the vectors

$$
\begin{aligned}
& \pm \beta_{i} \pm \beta_{j} \quad 1 \leq i, j \leq 7 \quad i \neq j \\
& \frac{1}{2}\left(\sum_{i=1}^{8} \varepsilon_{i} \beta_{i}\right) \quad \varepsilon_{i}= \pm 1, \quad \prod \varepsilon_{i}=1
\end{aligned}
$$

are all in the root system $\Phi$. We also have

$$
s_{8}\left(\beta_{i}\right)=\beta_{i}-2 \frac{\left\{\alpha_{8}, \beta_{i}\right\}}{\left\{\alpha_{8}, \alpha_{8}\right\}} \alpha_{8}=\beta_{i}+\frac{1}{2} \alpha_{8}
$$

for $1 \leq i \leq 8$. Thus

$$
s_{8}\left(\beta_{7}+\beta_{8}\right)=\frac{1}{2}\left(-\beta_{1}-\beta_{2}-\beta_{3}-\beta_{4}-\beta_{5}-\beta_{6}+\beta_{7}+\beta_{8}\right) \in \Phi
$$

Since $s_{8}^{2}=1$ it follows that $\beta_{7}+\beta_{8} \in \Phi$. We then see that $\pm \beta_{i} \pm \beta_{8} \in \Phi$ for all $i$ with $1 \leq i \leq 7$. Thus the set of vectors

$$
\begin{aligned}
& \pm \beta_{i} \pm \beta_{j} \\
& \frac{1}{2}\left(\sum_{i=1}^{8} \varepsilon_{i} \beta_{i}\right) \quad \varepsilon_{i}= \pm 1, \quad \Pi \varepsilon_{i}=1
\end{aligned}
$$

lies in $\Phi$. We shall show this is the full root system $\Phi$. In order to do so we must verify that this set is invariant under $s_{1}, \ldots, s_{8}$. It is clearly invariant under $s_{1}, \ldots, s_{7}$ since these fix $\beta_{8}$ and act by permutations together with an even number of sign changes on $\beta_{1}, \ldots, \beta_{7}$. Thus it is sufficient to verify that this set is invariant under $s_{8}$. Now we have

$$
s_{8}\left(\beta_{i}-\beta_{j}\right)=\beta_{i}-\beta_{j} \quad \text { for all } i \neq j
$$

Thus the set of vectors of form $\beta_{i}-\beta_{j}, i \neq j$, is invariant under $s_{8}$. Also

$$
s_{8}\left(\beta_{i}+\beta_{j}\right)=\beta_{i}+\beta_{j}+\alpha_{8} \quad \text { for all } i \neq j
$$

Thus $s_{8}$ transforms vectors of form $\beta_{i}+\beta_{j}, i \neq j$, into vectors $\frac{1}{2}\left(\sum \varepsilon_{i} \beta_{i}\right)$ with two $\varepsilon_{i}$ equal to 1 and six equal to -1 . Moreover all vectors $\frac{1}{2}\left(\sum \varepsilon_{i} \beta_{i}\right)$ with this property arise in this way. Similarly such vectors with $\operatorname{six} \varepsilon_{i}$ equal to 1 and two equal to -1 have the form $s_{8}\left(-\beta_{i}-\beta_{j}\right)$. Thus vectors of form $\beta_{i}+\beta_{j}$ or $-\beta_{i}-\beta_{j}$ with $i \neq j$ are transformed by $s_{8}$ into the given set, and so are vectors $\frac{1}{2} \sum \varepsilon_{i} \beta_{i}$ of type $(2,6)$ or $(6,2)$. The vectors of this form of type $(0,8)$ or $(8,0)$ are $\alpha_{8}$ and $-\alpha_{8}$, which are transformed into one another by $s_{8}$. It remains to show that $s_{8}$ transforms vectors $\frac{1}{2} \sum \varepsilon_{i} \beta_{i}$ of type $(4,4)$ into the given set. However, since such vectors have four positive signs and four negative signs they are orthogonal to $\alpha_{8}$, hence $s_{8}$ transforms each such
vector into itself. Thus the given set of vectors is invariant under $s_{1}, \ldots, s_{8}$ so is the full root system $\Phi$.

There are $4 \cdot\binom{8}{2}=112$ vectors of form $\pm \beta_{i} \pm \beta_{j}$ with $i \neq j$ and $2^{7}=128$ vectors of form $\frac{1}{2} \sum \varepsilon_{i} \beta_{i}$ with $\varepsilon_{i}= \pm 1$ and $\prod \varepsilon_{i}=1$. Thus the total number of roots is

$$
|\Phi|=112+128=240
$$

Finally we have $\operatorname{dim} L=8+|\Phi|=248$. Thus we have proved

Theorem 8.8 The simple Lie algebra of type $E_{8}$ has dimension 248.

We now turn to the simple Lie algebra of type $E_{7}$. Its Dynkin diagram is


Thus the vectors $\alpha_{2}, \alpha_{3}, \alpha_{4}, \alpha_{5}, \alpha_{6}, \alpha_{7}, \alpha_{8}$ considered above form a fundamental root system of type $E_{7}$. In order to obtain the full root system we must transform these vectors repeatedly by $s_{2}, \ldots, s_{8}$ until no new vectors are obtained. Now the vectors $\alpha_{2}, \ldots, \alpha_{8}$ are all orthogonal to $\beta_{1}-\beta_{8}$. Thus all their transforms by $s_{2}, \ldots, s_{8}$ will also be orthogonal to $\beta_{1}-\beta_{8}$. These transforms are contained in the set of roots of $E_{8}$ obtained above.

Now the roots of $E_{8}$ orthogonal to $\beta_{1}-\beta_{8}$ are:

$$
\begin{aligned}
& \pm \beta_{i} \pm \beta_{j}, \quad 2 \leq i, j \leq 7, \quad i \neq j \\
& \pm\left(\beta_{1}+\beta_{8}\right) \\
& \frac{1}{2} \sum \varepsilon_{i} \beta_{i}, \quad \varepsilon_{i}= \pm 1, \quad \prod \varepsilon_{i}=1, \quad \varepsilon_{1}=\varepsilon_{8}
\end{aligned}
$$

Thus the required root system of $E_{7}$ is contained in this set. We shall show it is the whole of this set.

We first consider the action of the subgroup of the Weyl group of $E_{7}$ generated by $s_{2}, s_{3}, s_{4}, s_{5}, s_{6}, s_{7}$. Elements of this subgroup fix $\beta_{1}$ and $\beta_{8}$ and act on $\beta_{2}, \beta_{3}, \beta_{4}, \beta_{5}, \beta_{6}, \beta_{7}$ by permutations combined with sign changes with an even number of negative signs. By applying elements of this subgroup to $\alpha_{2}, \ldots, \alpha_{8}$ we see that the vectors

$$
\begin{array}{ll} 
\pm \beta_{i} \pm \beta_{j}, & 2 \leq i, j \leq 7, \quad i \neq j \\
\frac{1}{2} \sum \varepsilon_{i} \beta_{i}, & \varepsilon_{i}= \pm 1, \quad \prod \varepsilon_{i}=1, \quad \varepsilon_{1}=\varepsilon_{8}
\end{array}
$$

are all roots of $E_{7}$. It remains to show that $\pm\left(\beta_{1}+\beta_{8}\right)$ are also roots of $E_{7}$. However,

$$
s_{8}\left(\beta_{1}+\beta_{8}\right)=\beta_{1}+\beta_{8}+\alpha_{8}=\frac{1}{2}\left(\beta_{1}-\beta_{2}-\beta_{3}-\beta_{4}-\beta_{5}-\beta_{6}-\beta_{7}+\beta_{8}\right)
$$

is a root of $E_{7}$, thus so is $\beta_{1}+\beta_{8}$ and $-\beta_{1}-\beta_{8}$.
There are $4\binom{6}{2}=60$ roots of form $\pm \beta_{i} \pm \beta_{j}, \quad 2 \leq i, j \leq 7, i \neq j$ and $2^{6}=$ 64 roots of form $\frac{1}{2} \sum \varepsilon_{i} \beta_{i}$ with $\varepsilon_{i}= \pm 1, \prod \varepsilon_{i}=1$ and $\varepsilon_{1}=\varepsilon_{8}$. Thus the number of roots of $E_{7}$ is given by

$$
|\Phi|=60+2+64=126 .
$$

Also we have

$$
\operatorname{dim} L=7+|\Phi|=133 .
$$

Thus we have shown:
Theorem 8.9 The simple Lie algebra of type $E_{7}$ has dimension 133.
Finally we consider the simple Lie algebra of type $E_{6}$. Its Dynkin diagram is


Thus the vectors $\alpha_{3}, \alpha_{4}, \alpha_{5}, \alpha_{6}, \alpha_{7}, \alpha_{8}$ considered above form a fundamental root system of type $E_{6}$. In order to obtain the full root system of $E_{6}$ we must transform these vectors successively by the fundamental reflections $s_{3}, s_{4}, s_{5}, s_{6}, s_{7}, s_{8}$.
Now the vectors $\alpha_{3}, \ldots, \alpha_{8}$ are all orthogonal to both $\beta_{1}-\beta_{8}$ and $\beta_{2}-\beta_{8}$. Thus the full root system of $E_{6}$ is orthogonal to $\beta_{1}-\beta_{8}$ and $\beta_{2}-\beta_{8}$.

Now the roots of $E_{8}$ orthogonal to both $\beta_{1}-\beta_{8}$ and $\beta_{2}-\beta_{8}$ are:

$$
\begin{aligned}
& \pm \beta_{i} \pm \beta_{j} \quad 3 \leq i, j \leq 7, \quad i \neq j \\
& \frac{1}{2}\left(\sum_{i=1}^{8} \varepsilon_{i} \beta_{i}\right) \quad \varepsilon_{i}= \pm 1, \quad \prod \varepsilon_{i}=1, \quad \varepsilon_{1}=\varepsilon_{2}=\varepsilon_{8} .
\end{aligned}
$$

Thus the required root system of $E_{6}$ is contained in this set. We shall show it is equal to this set of vectors.

Consider the action of the subgroup of the Weyl group of type $E_{6}$ generated by $s_{3}, s_{4}, s_{5}, s_{6}, s_{7}$. Elements of this subgroup fix $\beta_{1}, \beta_{2}, \beta_{8}$ and act on $\beta_{3}, \beta_{4}, \beta_{5}, \beta_{6}, \beta_{7}$ by permutations combined with sign changes with an even number of negative signs. By applying elements of this subgroup to $\alpha_{3}, \alpha_{4}, \alpha_{5}, \alpha_{6}, \alpha_{7}, \alpha_{8}$ we can obtain all vectors of form $\pm \beta_{i} \pm \beta_{j}$ with

Table 8.11 The simple Lie algebras

| $L$ |  | $\operatorname{dim} H$ | $\|\Phi\|$ | $\operatorname{dim} L$ |
| :--- | :---: | :---: | :---: | :---: |
| $A_{l} l \geq 1$ | $l$ | $l(l+1)$ | $l(l+2)$ |  |
| $B_{l}$ | $l \geq 2$ | $l$ | $2 l^{2}$ | $l(2 l+1)$ |
| $C_{l}$ | $l \geq 3$ | $l$ | $2 l^{2}$ | $l(2 l+1)$ |
| $D_{l}$ | $l \geq 4$ | $l$ | $2 l(l-1)$ | $l(2 l-1)$ |
| $E_{6}$ |  | 72 | 78 |  |
| $E_{7}$ | 7 | 126 | 133 |  |
| $E_{8}$ | 8 | 240 | 248 |  |
| $F_{4}$ | 4 | 48 | 52 |  |
| $G_{2}$ | 2 | 12 | 14 |  |

$3 \leq i, j \leq 7, i \neq j$, and (up to sign) all vectors of form $\frac{1}{2} \sum \varepsilon_{i} e_{i}$ with $\varepsilon_{i}= \pm 1$, $\prod \varepsilon_{i}=1, \varepsilon_{1}=\varepsilon_{2}=\varepsilon_{8}$. Hence the vectors in the above set are all roots of $E_{6}$. There are $\binom{5}{2} \cdot 4=40$ vectors of type $\pm \beta_{i} \pm \beta_{j}$ with $3 \leq i, j \leq 7, i \neq j$, and $2^{5}=32$ vectors of type $\frac{1}{2} \sum \varepsilon_{i} e_{i}$ with $\varepsilon_{i}= \pm 1, \prod \varepsilon_{i}=1$ and $\varepsilon_{1}=\varepsilon_{i}=\varepsilon_{8}$. Thus the total number of roots is

$$
|\Phi|=40+32=72
$$

and we have

$$
\operatorname{dim} L=6+|\Phi|=78
$$

Thus:

Theorem 8.10 The simple Lie algebra of type $E_{6}$ has dimension 78.
We have now determined the dimensions of all the simple Lie algebras. We summarise the information we have obtained in Table 8.11. In this table $L$ is a simple Lie algebra, $H$ is a Cartan subalgebra and $\Phi$ the system of roots of $L$ with respect to $H$.

### 8.8 Properties of long and short roots

Proposition 8.12 In the simple Lie algebras of types $A_{l}, D_{l}, E_{6}, E_{7}, E_{8}$ all the roots have the same length. In the Lie algebras of types $B_{l}, C_{l}, F_{4}, G_{2}$ there are two possible lengths of roots. These are called the long roots and short roots.

Proof. This is clear from the preceding results.

Proposition 8.13 (i) Let $\Phi$ be a root system of type $B_{l}$ with fundamental system


Then the long roots form a subsystem of type $D_{l}$ with fundamental system

and the short roots form a subsystem of type $\left(A_{1}\right)^{l}$ with fundamental system

$$
\begin{array}{cccc}
\mathrm{O} & \mathrm{O} & \cdots \cdot & \text { ○ } \\
\alpha_{1}+\cdots+\alpha_{l} & \alpha_{2}+\cdots+\alpha_{l} & & \alpha_{l-1}+\alpha_{l}
\end{array}
$$

(ii) Let $\Phi$ be a root system of type $C_{l}$ with fundamental system


Then the long roots form a subsystem of type $\left(A_{1}\right)^{l}$ with fundamental system

| ○ | O | $\cdots$ | O | O |
| :---: | :---: | :---: | :---: | :---: |
| $2 \alpha_{1}+\ldots+2 \alpha_{l-1}+\alpha_{l}$ | $2 \alpha_{2}+\ldots+2 \alpha_{l-1}+\alpha_{l}$ |  | $2 \alpha_{l-1}+\alpha_{l}$ | $\alpha_{l}$ |

and the short roots form a subsystem of type $D_{l}$ with fundamental system

(iii) Let $\Phi$ be a root system of type $F_{4}$ with fundamental system


Then the long roots form a subsystem of type $D_{4}$ with fundamental system

and the short roots form a subsystem of type $D_{4}$ with fundamental system

(iv) Let $\Phi$ be a root system of type $G_{2}$ with fundamental system


Then the long roots form a subsystem of type $A_{2}$ with fundamental system

$$
\stackrel{O}{\alpha_{1}}-\alpha_{1}+3 \alpha_{2}
$$

and the short roots form a subsystem of type $A_{2}$ with fundamental system

$$
\stackrel{\mathrm{O}}{\alpha_{1}+\alpha_{2}} \alpha_{2}^{-}
$$

Proof. (i) We saw in Section 8.3 that the roots of type $B_{l}$ have form $\pm \beta_{i} \pm \beta_{j}, i \neq j$, and $\pm \beta_{i}$. The former are the long roots and the latter the short roots. The long roots form a system of type $D_{l}$ with fundamental system $\beta_{1}-\beta_{2}, \beta_{2}-\beta_{3}, \ldots, \beta_{l-1}-\beta_{l}, \beta_{l-1}+\beta_{l}$. These are $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{l-1}$, $\alpha_{l-1}+2 \alpha_{l}$ respectively. The short roots form a system of type $\left(A_{1}\right)^{l}$ with fundamental system $\beta_{1}, \ldots, \beta_{l}$. These are $\alpha_{1}+\cdots+\alpha_{l}, \alpha_{2}+\cdots+\alpha_{l}, \ldots, \alpha_{l}$ respectively.
(ii) We saw in Section 8.4 that the roots of type $C_{l}$ have form $\pm \beta_{i} \pm \beta_{j}, i \neq j$, and $\pm 2 \beta_{i}$. The former are the short roots and the latter the long roots. Thus the short roots form a subsystem of type $D_{l}$ and the long roots a subsystem of type $\left(A_{1}\right)^{l}$.
(iii) We saw in Section 8.6 that the roots of type $F_{4}$ have form

$$
\begin{aligned}
& \pm \beta_{i} \pm \beta_{j} \quad i \neq j \\
& \pm \beta_{i} \\
& \frac{1}{2}\left( \pm \beta_{1} \pm \beta_{2} \pm \beta_{3} \pm \beta_{4}\right) .
\end{aligned}
$$

Roots of the first type are long and those of the second and third types are short. The long roots form a subsystem of type $D_{4}$ with fundamental system $\beta_{4}-\beta_{1}, \beta_{1}-\beta_{2}, \beta_{2}-\beta_{3}, \beta_{2}+\beta_{3}$. These are $\alpha_{2}+2 \alpha_{3}+2 \alpha_{4}, \alpha_{1}, \alpha_{2}$,
$\alpha_{2}+2 \alpha_{3}$ respectively. The short roots also form a subsystem of type $D_{4}$, with fundamental system $\beta_{1}, \beta_{2}, \beta_{3}, \frac{1}{2}\left(-\beta_{1}-\beta_{2}-\beta_{3}+\beta_{4}\right)$. These are $\alpha_{1}+\alpha_{2}+\alpha_{3}, \alpha_{2}+\alpha_{3}, \alpha_{3}, \alpha_{4}$ respectively.
(iv) The long and short roots of type $G_{2}$ are evident from Section 8.5.

Let $\Pi=\left\{\alpha_{i}\right\}$ be a fundamental system of roots in a simple Lie algebra whose Dynkin diagram has a double or triple edge, and let $\Phi$ be the root system with fundamental system $\Pi$. Consider the simple Lie algebra whose Dynkin diagram is obtained from that above by reversing the direction of the arrow. Let $\Pi^{\mathrm{v}}=\left\{\alpha_{i}^{\mathrm{v}}\right\}$ be the corresponding fundamental system, labelled as before, and $\Phi^{\mathrm{v}}$ be the root system with fundamental system $\Pi^{\mathrm{v}}$.

System $\Phi^{v}$ is called the dual root system of $\Phi$. The possible types of $\Phi, \Phi^{\mathrm{v}}$ are as shown.

| $\underline{\Phi}$ | $\underline{\Phi}^{\mathrm{v}}$ |
| :--- | :--- |
| $B_{l}$ | $C_{l}$ |
| $C_{l}$ | $B_{l}$ |
| $F_{4}$ | $F_{4}$ |
| $G_{2}$ | $G_{2}$ |

We note that $\alpha_{i}$ is a short root in $\Pi$ if and only if $\alpha_{i}^{v}$ is a long root in $\Pi^{v}$.
We suppose as usual that we have symmetric scalar products $\{$,$\} on \mathbb{R} \Pi$ and $\mathbb{R} \Pi^{\mathrm{v}}$ such that

$$
2 \frac{\left\{\alpha_{i}, \alpha_{j}\right\}}{\left\{\alpha_{i}, \alpha_{i}\right\}}=A_{i j}, \quad 2 \frac{\left\{\alpha_{i}^{\mathrm{v}}, \alpha_{j}^{\mathrm{v}}\right\}}{\left\{\alpha_{i}^{\mathrm{v}}, \alpha_{i}^{\mathrm{v}}\right\}}=A_{i j}^{\mathrm{v}}
$$

for all $i, j$, where $A=\left(A_{i j}\right), A^{\mathrm{v}}=\left(A_{i j}^{\mathrm{v}}\right)$ are the Cartan matrices of $\Phi, \Phi^{\mathrm{v}}$ respectively.

We consider the free abelian groups $\mathbb{Z} \Pi, \mathbb{Z} \Pi^{v}$ generated by $\Pi, \Pi^{v}$. We define a homomorphism

$$
\theta: \mathbb{Z} \Pi^{\mathrm{v}} \rightarrow \mathbb{Z} \Pi
$$

by

$$
\theta\left(\alpha_{i}^{v}\right)= \begin{cases}p \alpha_{i} & \text { if } \alpha_{i} \text { is a short root } \\ \alpha_{i} & \text { if } \alpha_{i} \text { is a long root }\end{cases}
$$

where $p$ is the ratio of the squared lengths of the long and short roots. (Thus $p=2$ in types $B_{l}, C_{l}, F_{4}$ and $p=3$ in type $G_{2}$.)

Lemma $8.142 \frac{\left\{\theta\left(\alpha_{i}^{v}\right), \theta\left(\alpha_{j}^{v}\right)\right\}}{\left\{\theta\left(\alpha_{i}^{v}\right), \theta\left(\alpha_{i}^{v}\right)\right\}}=A_{i j}^{\mathrm{v}}$ for all $i$, $j$.
Proof. We write $\theta\left(\alpha_{i}^{\mathrm{v}}\right)=\xi_{i} \alpha_{i}$ where $\xi_{i}=p$ if $\alpha_{i}$ is short and $\xi_{i}=1$ if $\alpha_{i}$ is long. Then

$$
2 \frac{\left\{\theta\left(\alpha_{i}^{v}\right), \theta\left(\alpha_{j}^{v}\right)\right\}}{\left\{\theta\left(\alpha_{i}^{v}\right), \theta\left(\alpha_{i}^{v}\right)\right\}}=\xi_{i}^{-1} \xi_{j} 2 \frac{\left\{\alpha_{i}, \alpha_{j}\right\}}{\left\{\alpha_{i}, \alpha_{i}\right\}}=\xi_{i}^{-1} \xi_{j} A_{i j}
$$

However, $\xi_{i}^{-1} \xi_{j} A_{i j}=A_{i j}^{\mathrm{v}}$. This follows from the following observations.
If $A_{i j} \neq 0$ and $\alpha_{i}, \alpha_{j}$ have the same length then we have $\xi_{i}=\xi_{j}$ and $A_{i j}=$ $A_{i j}^{\mathrm{v}}=-1$.

If $A_{i j} \neq 0, \alpha_{i}$ is long, $\alpha_{j}$ is short then $\xi_{i}=1, \xi_{j}=p, A_{i j}=-1$ and $A_{i j}^{\mathrm{v}}=-p$.
If $A_{i j} \neq 0, \alpha_{i}$ is short, $\alpha_{j}$ is long then $\xi_{i}=p, \xi_{j}=1, A_{i j}=-p, A_{i j}^{\mathrm{v}}=-1$. Thus in all cases we have $\xi_{i}^{-1} \xi_{j} A_{i j}=A_{i j}^{\mathrm{v}}$.

Lemma 8.15 The diagram

commutes.

Proof. On the one hand $\theta s_{i}^{\mathrm{v}}\left(\alpha_{j}^{\mathrm{v}}\right)=\theta\left(\alpha_{j}^{\mathrm{v}}-A_{i j}^{\mathrm{v}} \alpha_{i}^{\mathrm{v}}\right)=\theta\left(\alpha_{j}^{\mathrm{v}}\right)-A_{i j}^{\mathrm{v}} \theta\left(\alpha_{i}^{\mathrm{v}}\right)$. On the other hand

$$
\begin{aligned}
s_{i} \theta\left(\alpha_{j}^{\mathrm{v}}\right) & =s_{i}\left(\xi_{j} \alpha_{j}\right)=\xi_{j} s_{i}\left(\alpha_{j}\right)=\xi_{j}\left(\alpha_{j}-A_{i j} \alpha_{i}\right) \\
& =\xi_{j} \alpha_{j}-\xi_{i}^{-1} \xi_{j} A_{i j}\left(\xi_{i} \alpha_{i}\right)=\theta\left(\alpha_{j}^{\mathrm{v}}\right)-A_{i j}^{\mathrm{v}} \theta\left(\alpha_{i}^{\mathrm{v}}\right)
\end{aligned}
$$

Thus $\theta s_{i}^{\mathrm{v}}=s_{i} \theta$.
Let $W, W^{\mathrm{v}}$ be the Weyl groups of $\Phi, \Phi^{\mathrm{v}}$. There is a natural isomorphism $W \cong W^{\mathrm{v}}$ under which $s_{i}$ corresponds to $s_{i}^{\mathrm{v}}$, since the root lengths are irrelevant as far as the structure of the Weyl group is concerned. We shall use this isomorphism to identify $W^{v}$ with $W$. Then Lemma 8.15 shows that $\theta w=w \theta$ for all $w \in W$.

Proposition 8.16 Given $\alpha^{\mathrm{v}} \in \Phi^{v}$ there is a unique $\alpha \in \Phi$ such that

$$
\theta\left(\alpha^{\mathrm{v}}\right)= \begin{cases}p \alpha & \text { if } \alpha \text { is short } \\ \alpha & \text { if } \alpha \text { is long }\end{cases}
$$

Proof. We have $\alpha^{\mathrm{v}}=w\left(\alpha_{i}^{\mathrm{v}}\right)$ for some $w \in W, \alpha_{i}^{\mathrm{v}} \in \Pi^{\mathrm{v}}$. Thus

$$
\theta\left(\alpha^{v}\right)=\theta w\left(\alpha_{i}^{\mathrm{v}}\right)=w \theta\left(\alpha_{i}^{\mathrm{v}}\right)= \begin{cases}p w\left(\alpha_{i}\right) & \text { if } \alpha_{i} \text { is short } \\ w\left(\alpha_{i}\right) & \text { if } \alpha_{i} \text { is long }\end{cases}
$$

Let $\alpha=w\left(\alpha_{i}\right) \in \Phi$. Then $\alpha$ is uniquely determined since

$$
w\left(\alpha_{i}^{\mathrm{v}}\right)=w^{\prime}\left(\alpha_{j}^{\mathrm{v}}\right) \Rightarrow w^{\prime-1} w\left(\alpha_{i}^{\mathrm{v}}\right)=\alpha_{j}^{\mathrm{v}} \Rightarrow w^{\prime-1} w\left(\alpha_{i}\right)=\alpha_{j} \Rightarrow w\left(\alpha_{i}\right)=w^{\prime}\left(\alpha_{j}\right)
$$

Thus $\theta\left(\alpha^{v}\right)= \begin{cases}p \alpha & \text { if } \alpha \text { is short } \\ \alpha & \text { if } \alpha \text { is long. }\end{cases}$
This proposition determines a bijection $\Phi^{v} \rightarrow \Phi$ under which $\alpha^{v} \rightarrow \alpha . \alpha^{v}$ is called the dual root of $\alpha . \alpha^{\mathrm{v}}$ is long if and only if $\alpha$ is short.

Proposition 8.17 Let $\alpha \in \Phi$ satisfy $\alpha=\sum n_{i} \alpha_{i}$. Then $\alpha$ is a long root if and only if $p$ divides $n_{i}$ for all $i$ for which $\alpha_{i}$ is a short root.

Proof. Suppose $\alpha$ is long. Then $\theta\left(\alpha^{v}\right)=\alpha$ and so

$$
\alpha^{\mathrm{v}}=\sum_{\alpha_{i} \text { long }} n_{i} \alpha_{i}^{\mathrm{v}}+\sum_{\alpha_{i} \text { short }} n_{i} p^{-1} \alpha_{i}^{\mathrm{v}}
$$

Since $\alpha^{\mathrm{v}} \in \sum_{i} \mathbb{Z} \alpha_{i}^{\mathrm{v}}$ we deduce that $p$ divides $n_{i}$ whenever $\alpha_{i}$ is short.
Now suppose conversely that $p$ divides $n_{i}$ for all $i$ for which $\alpha_{i}$ is short. Let $n_{i}=p m_{i}$ for such $i$. Suppose if possible that $\alpha$ is short. Then $\theta\left(\alpha^{v}\right)=p \alpha$. Thus

$$
\begin{aligned}
\alpha^{\mathrm{v}} & =\sum_{\alpha_{i} \text { long }} p n_{i} \alpha_{i}^{\mathrm{v}}+\sum_{\alpha_{i} \text { short }} n_{i} \alpha_{i}^{\mathrm{v}} \\
& =p\left(\sum_{\alpha_{i} \text { long }} n_{i} \alpha_{i}^{\mathrm{v}}+\sum_{\alpha_{i} \text { short }} m_{i} \alpha_{i}^{\mathrm{v}}\right) .
\end{aligned}
$$

This gives $\alpha^{\mathrm{v}} \in \sum p \mathbb{Z} \alpha_{i}^{\mathrm{v}}$ which is impossible. Thus $\alpha$ is a long root.

Proposition 8.18 (i) The abelian group generated by the short roots in $\Phi$ is $\sum_{i=1}^{l} \mathbb{Z} \alpha_{i}$.
(ii) The abelian group generated by the long roots in $\Phi$ is $\sum_{\alpha_{i} \text { long }} \mathbb{Z} \alpha_{i}+$ $\sum_{\alpha_{i} \text { short }} p \mathbb{Z} \alpha_{i}$.

Proof. By considering the fundamental system of the subsystem of short roots described in Proposition 8.13 it is clear that the abelian group generated by the short roots contains $\alpha_{1}, \ldots, \alpha_{l}$ so is $\sum \mathbb{Z} \alpha_{i}$.

The abelian group generated by the long roots lies in $\sum_{\alpha_{i} \text { long }} \mathbb{Z} \alpha_{i}+$ $\sum_{\alpha_{i} \text { short }} p \mathbb{Z} \alpha_{i}$ by Proposition 8.17. However, this group contains $\alpha_{i}$ for $\alpha_{i}$ long and $p \alpha_{i}$ for $\alpha_{i}$ short, again by Proposition 8.17 , so must be $\sum_{\alpha_{i} \text { long }} \mathbb{Z} \alpha_{i}+$ $\sum_{\alpha_{i} \text { short }} p \mathbb{Z} \alpha_{i}$.

## 9

## Some universal constructions

### 9.1 The universal enveloping algebra

Let $L$ be a Lie algebra over $\mathbb{C}$. We shall show in this section how to construct an associative algebra $\mathfrak{H}(L)$, the universal enveloping algebra of $L$, such that the representation theory of $\mathfrak{H}(L)$ is the same as the representation theory of the Lie algebra $L$. Even if $L$ is finite dimensional its enveloping algebra $\mathfrak{H}(L)$ will be infinite dimensional.

We begin by forming the tensor powers of $L$. We define $T^{0}$ to be the 1-dimensional vector space $\mathbb{C} 1, T^{1}=L, T^{2}=L \otimes_{\mathbb{C}} L$ and, in general,

$$
T^{n}=\mathrm{L} \otimes_{\mathbb{C}} L \otimes_{\mathbb{C}} \cdots \otimes_{\mathbb{C}} L \quad(n \text { factors })
$$

$T^{n}$ is a vector space over $\mathbb{C}$ of dimension $(\operatorname{dim} L)^{n}$.
We next form the tensor algebra $T=T(L)$ of $L$. We define $T$ as the direct sum of vector spaces

$$
T=T^{0} \oplus T^{1} \oplus T^{2} \oplus \cdots
$$

Thus elements of $T$ are finite sums of elements, each of which lies in some $T^{n}$. We may define a bilinear map

$$
T^{m} \times T^{n} \rightarrow T^{m+n}
$$

satisfying

$$
\left(x_{1} \otimes \cdots \otimes x_{m}\right) \cdot\left(y_{1} \otimes \cdots \otimes y_{n}\right)=x_{1} \otimes \cdots \otimes x_{m} \otimes y_{1} \otimes \cdots \otimes y_{n}
$$

for $x_{i}, y_{j} \in L$ and then extend this map by linearity to give a multiplication map

$$
T \times T \rightarrow T
$$

In this way $T$ becomes an associative algebra called the tensor algebra of $L$. The element $1 \in T^{0}$ is the identity element of $T$.

Let $J$ be the 2-sided ideal of $T$ generated by all elements of the form

$$
x \otimes y-y \otimes x-[x y] \quad \text { for } x, y \in L
$$

$J$ is in particular a subspace of $T$. Let $\mathfrak{H}(L)=T / J$. Then $\mathfrak{H}(L)$ is an associative algebra over $\mathbb{C}$ called the universal enveloping algebra of $L$.

Example 9.1 Let $L$ be an $n$-dimensional abelian Lie algebra over $\mathbb{C}$. Then $L$ has basis $x_{1}, \ldots, x_{n}$ and we have $\left[x_{i} x_{j}\right]=0$ for all $i, j$. Thus $J$ is the 2 -sided ideal of $T$ generated by all elements of the form $x \otimes y-y \otimes x$ for $x, y \in L$. Thus $\mathfrak{H}(L)$ is a commutative algebra, and is generated as an algebra by the identity 1 and the elements $x_{1}, \ldots, x_{n}$. In fact $\mathfrak{H}(L)$ is isomorphic to the polynomial algebra $\mathbb{C}\left[x_{1}, \ldots, x_{n}\right]$.

In general we have linear maps

$$
L \rightarrow T^{1} \rightarrow T \rightarrow \mathfrak{U}(L)
$$

and we denote by $\sigma: \mathrm{L} \rightarrow \mathfrak{U}(L)$ the composite linear map. We now show that $\mathfrak{H}(L)$ has a certain universal property which justifies its name.

Proposition 9.2 Let $A$ be any associative algebra with 1 over $\mathbb{C}$ and $[A]$ the corresponding Lie algebra. Then given any Lie algebra homomorphism $\theta: L \rightarrow[A]$ there exists a unique associative algebra homomorphism $\phi: \mathfrak{U}(L) \rightarrow A$ such that $\phi \circ \sigma=\theta$.

Note Associative algebra homomorphisms will be understood to be homomorphisms of associative algebras with identity in this chapter. Thus the homomorphism will map identity to identity.

Proof. We first observe that the linear map $\theta: L \rightarrow A$ can be extended to an associative algebra homomorphism from $T$ to $A$. If $x_{i}, i \in I$, are a basis for $L$ then the set of all monomials $x_{i_{1}} \ldots x_{i_{r}}$ for $i_{1}, \ldots, i_{r} \in I$ form a basis for $T$. The case $r=0$ gives the identity element. The map

$$
x_{i_{1}} \ldots x_{i_{r}} \rightarrow \theta\left(x_{i_{1}}\right) \ldots \theta\left(x_{i_{r}}\right)
$$

can then be extended by linearity to give an associative algebra homomorphism from $T$ to $A$. Let this map be $\theta^{\prime}: T \rightarrow A$. Let $x, y \in L$. Then we have

$$
\begin{aligned}
\theta^{\prime}(x & \otimes y-y \otimes x-[x y]) \\
& =\theta(x) \theta(y)-\theta(y) \theta(x)-\theta[x y] \\
& =[\theta(x), \theta(y)]-\theta[x y]=0
\end{aligned}
$$

since $\theta: L \rightarrow[A]$ is a Lie algebra homomorphism. Thus all the generators of the 2 -sided ideal $J$ of $T$ lie in the kernel of $\theta^{\prime}$. Since the kernel is a 2 -sided ideal, $J$ lies in the kernel of $\theta^{\prime}$. This shows there is an induced homomorphism

$$
\phi: T / J \rightarrow A
$$

such that the diagram

commutes. When we restrict the domain to $T^{1}$ we deduce that $\phi \circ \sigma=\theta$. This proves the existence of a homomorphism $\phi: \mathfrak{U}(L) \rightarrow A$ of the required type.

We now prove the uniqueness of $\phi$. Let $\phi^{\prime}: \mathfrak{U}(L) \rightarrow A$ be another such homomorphism. Now $T$ is generated by $T^{1}$ as an associative algebra with 1. Thus its factor algebra $\mathfrak{H}(L)$ is generated by $\sigma(L)$, which is the image of $T^{1}$ in $\mathfrak{H}(L)$. Let $x \in L$. Then

$$
\phi^{\prime}(\sigma(x))=\theta(x)=\phi(\sigma(x))
$$

Thus $\phi, \phi^{\prime}$ agree on $\sigma(x)$ for all $x \in L$. Since $\sigma(L)$ generates $\mathfrak{l}(L)$ it follows that $\phi, \phi^{\prime}$ agree on $\mathfrak{H}(L)$, so $\phi^{\prime}=\phi$.

Using this universal property we can relate representations of the Lie algebra $L$ to representations of the associative algebra $\mathfrak{l}(L)$. If $V$ is a vector space over $\mathbb{C}$ the set End $V$ of all linear maps of $V$ into itself forms an associative algebra with 1 , and the corresponding Lie algebra is [End $V$ ]. A representation of $L$ is a Lie algebra homomorphism $L \rightarrow[$ End $V]$ and a representation of $\mathfrak{H}(L)$ is an associative algebra homomorphism $\mathfrak{H}(L) \rightarrow$ End $V$.

Proposition 9.3 There is a bijective correspondence between representations $\theta: L \rightarrow[$ End $V]$ and representations $\phi: \mathfrak{H}(L) \rightarrow$ End $V$. Corresponding representations are related by the condition

$$
\phi(\sigma(x))=\theta(x) \quad \text { for all } x \in L
$$

Proof. Let $\theta: L \rightarrow[$ End $V]$ be a representation of $L$. Then by Proposition 9.2 there exists a unique associative algebra homomorphism $\phi: \mathfrak{U}(L) \rightarrow$ End $V$ such that $\phi \circ \sigma=\theta$.

Conversely, given an associative algebra homomorphism $\phi: \mathfrak{H}(L) \rightarrow$ End $V$ we wish to define a corresponding Lie algebra homomorphism
$\theta: L \rightarrow[$ End $V]$. Now we have a linear map $\sigma: L \rightarrow \mathfrak{H}(L)$. Since $\mathfrak{H}(L)=$ $T / J$ and $x \otimes y-y \otimes x-[x y] \in J$ for all $x, y \in L$ we see that

$$
\sigma(x) \sigma(y)-\sigma(y) \sigma(x)-\sigma[x y]=0
$$

for $x, y \in L$. This gives

$$
[\sigma(x), \sigma(y)]=\sigma[x y]
$$

and so $\sigma: L \rightarrow[\mathfrak{H}(L)]$ is a Lie algebra homomorphism. We now define $\theta: L \rightarrow$ [End $V$ ] by $\theta=\phi \circ \sigma$. Then $\theta$ is a Lie homomorphism of the required type.

It is clear from the definitions that the maps $\theta \rightarrow \phi$ and $\phi \rightarrow \theta$ are inverse to one another.

We shall find this result very useful in the subsequent development, when we shall obtain information about representations of finite dimensional Lie algebras by considering the representation theory of the corresponding universal enveloping algebra.

### 9.2 The Poincaré-Birkhoff-Witt basis theorem

We shall now describe how to obtain a basis for the universal enveloping algebra $\mathfrak{l}(L)$.

Theorem 9.4 (Poincaré-Birkhoff-Witt). Let L be a Lie algebra with basis $\left\{x_{i} ; i \in I\right\}$. Let $<$ be a total order on the index set I. Let $\sigma: L \rightarrow \mathfrak{U}(L)$ be the natural linear map from $L$ into its enveloping algebra. Let $\sigma\left(x_{i}\right)=y_{i}$. Then the elements

$$
y_{i_{1}}^{r_{1}} \ldots y_{i_{n}}^{r_{n}}
$$

for all $n \geq 0$, all $r_{i} \geq 0$, and all $i_{1}, \ldots, i_{n} \in I$ with $i_{1}<i_{2}<\cdots<i_{n}$ form a basis for $\mathfrak{H}(L)$.

Proof. (a) We first show that the above elements $y_{i_{1}}^{r_{1}} \ldots y_{i_{n}}^{r_{n}}$ span $\mathfrak{H}(L)$. We know that the elements of form $x_{j_{1}} \otimes \ldots \otimes x_{j_{k}}$ for all $k$ and all $j_{1}, \ldots, j_{k} \in I$ span $T$. By applying the natural homomorphism $T \rightarrow \mathfrak{H}(L)$ it follows that the elements of form $y_{j_{1}} \ldots y_{j_{k}}$ span $\mathfrak{U}(L)$. It is therefore sufficient to show that every product $y_{j_{1}} \ldots y_{j_{k}}$ is a linear combination of the given elements of form $y_{i_{1}}^{r_{1}} \ldots y_{i_{n}}^{r_{n}}$. We shall prove this by induction on $k$. It is obvious if $k=1$. For arbitrary $k$ it is clear when $j_{1} \leq \cdots \leq j_{k}$. If this is not so we may use relations of form

$$
y_{i} y_{j}=y_{j} y_{i}+\sigma\left[x_{i} x_{j}\right] .
$$

We note that $\left[x_{i} x_{j}\right]$ is a linear combination of elements $x_{t}$ for $t \in I$ and so $\sigma\left[x_{i} x_{j}\right]$ is a linear combination of $y_{t}$ for $t \in I$. Thus we may interchange the order of two consecutive terms $y_{i}, y_{j}$ in a monomial of degree $k$ provided we introduce a certain linear combination of monomials of degree less than $k$. By performing such interchanges a finite number of times we may express the terms $y_{i}$ in the monomial with the $i$ in the given order $<$ on $I$. Thus

$$
\begin{gathered}
y_{j_{1}} \ldots y_{j_{k}}=y_{i_{1}}^{r_{1}} \ldots y_{i_{n}}^{r_{n}}+\text { a linear combination of monomials of degree } \\
\text { less than } k
\end{gathered}
$$

where $r_{1}+\cdots+r_{n}=k$ and $i_{1}<\cdots<i_{n}$. By induction we may assume that all monomials of degree less than $k$ are expressible as linear combinations of monomials with terms in the given order $<$. The required result then follows.
(b) We now show that the given monomials of form $y_{i_{1}}^{r_{1}} \ldots y_{i_{n}}^{r_{n}}$ are linearly independent. This is not so easy to see, and we shall prove it by an indirect argument. We introduce the polynomial ring $R=\mathbb{C}\left[z_{i} ; i \in I\right]$ and shall make use of the following lemma.

Lemma 9.5 There exists a linear map $\theta: T \rightarrow R$ satisfying the conditions

$$
\begin{aligned}
& \theta\left(x_{i_{1}} \otimes \cdots \otimes x_{i_{n}}\right)=z_{i_{1}} \ldots z_{i_{n}} \quad \text { if } i_{1} \leq \cdots \leq i_{n} \\
& \theta\left(x_{i_{1}} \otimes \cdots \otimes x_{i_{k}} \otimes x_{i_{k+1}} \otimes \cdots \otimes x_{i_{n}}-x_{i_{1}} \otimes \cdots \otimes x_{i_{k+1}} \otimes x_{i_{k}} \otimes \cdots \otimes x_{i_{n}}\right) \\
& \quad=\theta\left(x_{i_{1}} \otimes \cdots \otimes\left[x_{i_{k}} x_{i_{k+1}}\right] \otimes \cdots \otimes x_{i_{n}}\right) \quad \text { for all } i_{1}, \ldots, i_{n} \text { and all } k \\
& \quad \text { with } 1 \leq k<n .
\end{aligned}
$$

Proof. We define the index of the monomial $x_{i_{1}} \otimes \cdots \otimes x_{i_{n}}$ to be the number of pairs ( $r, s$ ) with $1 \leq r<s \leq n$ satisfying $i_{r}>i_{s}$. Thus the monomials of index 0 are those whose terms appear in their natural order. Let $T^{n, j}$ be the subspace of $T^{n}$ spanned by all monomials $x_{i_{1}} \otimes \cdots \otimes x_{i_{n}}$ of index at most $j$. Thus

$$
T^{n, 0} \subset T^{n, 1} \subset \cdots \subset T^{n}
$$

We define $\theta: T^{0} \rightarrow R$ by $\theta(1)=1$. Suppose inductively that $\theta: T^{0} \oplus \cdots \oplus$ $T^{n-1} \rightarrow R$ has already been defined satisfying the required conditions. We shall show that $\theta$ can be extended to $\theta: T^{0} \oplus \cdots \oplus T^{n} \rightarrow R$. We define $\theta: T^{n, 0} \rightarrow R$ by

$$
\theta\left(x_{i_{1}} \otimes \cdots \otimes x_{i_{n}}\right)=z_{i_{1}} \ldots z_{i_{n}}
$$

if the monomial $x_{i_{1}} \otimes \cdots \otimes x_{i_{n}}$ has index 0 . We suppose $\theta: T^{n, i} \rightarrow R$ has already been defined, thus giving a linear map from $T^{0} \oplus \cdots \oplus T^{n-1} \oplus T^{n, i}$ to
$R$ satisfying the required conditions. We wish to define $\theta: T^{n, i+1} \rightarrow R$. Thus suppose the monomial $x_{i_{1}} \otimes \cdots \otimes x_{i_{n}}$ has index $i+1$. Then there exists $k$ with $1 \leq k<n$ such that $x_{i_{1}} \otimes \cdots \otimes x_{i_{k-1}} \otimes x_{i_{k+1}} \otimes x_{i_{k}} \otimes x_{i_{k+2}} \otimes \cdots \otimes x_{i_{n}}$ has index $i$. We then wish to define $\theta\left(x_{i_{1}} \otimes \cdots \otimes x_{i_{n}}\right)$ by the formula

$$
\begin{aligned}
\theta\left(x_{i_{1}}\right. & \left.\otimes \cdots \otimes x_{i_{k}} \otimes x_{i_{k+1}} \otimes \cdots \otimes x_{i_{n}}\right) \\
= & \theta\left(x_{i}, \otimes \cdots \otimes x_{i_{k+1}} \otimes x_{i_{k}} \otimes \cdots \otimes x_{i_{n}}\right) \\
& +\theta\left(x_{i_{1}} \otimes \cdots \otimes\left[x_{i_{k}} x_{i_{k+1}}\right] \otimes \cdots \otimes x_{i_{n}}\right)
\end{aligned}
$$

noting that the terms on the right-hand side have already been defined. However, there may be more than one possible choice of $k$ and we must check that if we choose a different one the linear map $\theta: T^{n, i+1} \rightarrow R$ will still be the same. So suppose $k^{\prime}$ also satisfies $1 \leq k^{\prime}<n$. We may without loss of generality assume that $k<k^{\prime}$.

We suppose first that $k+1<k^{\prime}$. Let $x_{i_{k}}=a, x_{i_{k+1}}=b, x_{i_{k^{\prime}}}=c, x_{i_{k^{\prime}+1}}=d$. Then the definition using the integer $k$ gives

$$
\begin{aligned}
\theta(\cdots \otimes & a \otimes b \otimes \cdots \otimes c \otimes d \otimes \cdots) \\
= & \theta(\cdots \otimes b \otimes a \otimes \cdots \otimes c \otimes d \otimes \cdots) \\
& +\theta(\cdots \otimes[a b] \otimes \cdots \otimes c \otimes d \otimes \cdots) \\
= & \theta(\cdots \otimes b \otimes a \otimes \cdots \otimes d \otimes c \otimes \cdots) \\
& +\theta(\cdots \otimes b \otimes a \otimes \cdots \otimes[c d] \otimes \cdots) \\
& +\theta(\cdots \otimes[a b] \otimes \cdots \otimes d \otimes c \otimes \cdots) \\
& +\theta(\cdots \otimes[a b] \otimes \cdots \otimes[c d] \otimes \cdots)
\end{aligned}
$$

using the inductive assumptions.
The second definition using the integer $k^{\prime}$ gives

$$
\begin{aligned}
\theta(\cdots \otimes & a \otimes b \otimes \cdots \otimes c \otimes d \otimes \cdots) \\
= & \theta(\cdots \otimes a \otimes b \otimes \cdots \otimes d \otimes c \otimes \cdots) \\
& +\theta(\cdots \otimes a \otimes b \otimes \cdots \otimes[c d] \otimes \cdots) \\
= & \theta(\cdots \otimes b \otimes a \otimes \cdots \otimes d \otimes c \otimes \cdots) \\
& +\theta(\cdots \otimes[a b] \otimes \cdots \otimes d \otimes c \otimes \cdots) \\
& +\theta(\cdots \otimes b \otimes a \otimes \cdots \otimes[c d] \otimes \cdots) \\
& +\theta(\cdots \otimes[a b] \otimes \cdots \otimes[c d] \otimes \cdots)
\end{aligned}
$$

using the inductive assumptions. These two expressions using integers $k, k^{\prime}$ are the same.

Now suppose that $k^{\prime}=k+1$. Let $x_{i_{k}}=a, x_{i_{k+1}}=b, x_{i_{k+2}}=c$. We compare the two ways of calculating $\theta(\cdots \otimes a \otimes b \otimes c \otimes \cdots)$. The first method, using the integer $k$, gives

$$
\begin{aligned}
\theta(\cdots \otimes & a \otimes b \otimes c \otimes \cdots) \\
= & \theta(\cdots \otimes b \otimes a \otimes c \otimes \cdots)+\theta(\cdots \otimes[a b] \otimes c \otimes \cdots) \\
= & \theta(\cdots \otimes b \otimes c \otimes a \otimes \cdots)+\theta(\cdots \otimes b \otimes[a c] \otimes \cdots) \\
& +\theta(\cdots \otimes c \otimes[a b] \otimes \cdots)+\theta(\cdots \otimes[[a b] c] \otimes \cdots) \\
= & \theta(\cdots \otimes c \otimes b \otimes a \otimes \cdots)+\theta(\cdots \otimes[b c] \otimes a \otimes \cdots) \\
& +\theta(\cdots \otimes b \otimes[a c] \otimes \cdots)+\theta(\cdots \otimes c \otimes[a b] \otimes \cdots) \\
& +\theta(\cdots \otimes[[a b] c] \otimes \cdots) \\
= & \theta(\cdots \otimes c \otimes b \otimes a \otimes \cdots)+\theta(\cdots \otimes a \otimes[b c] \otimes \cdots) \\
& +\theta(\cdots \otimes b \otimes[a c] \otimes \cdots)+\theta(\cdots \otimes c \otimes[a b] \otimes \cdots) \\
& +\theta(\cdots \otimes[[a b] c] \otimes \cdots)+\theta(\cdots \otimes[[b c] a] \otimes \cdots)
\end{aligned}
$$

using the inductive assumptions. The second method, using the integer $k^{\prime}=$ $k+1$, gives

$$
\begin{aligned}
\theta(\cdots \otimes & a \otimes b \otimes c \otimes \cdots) \\
= & \theta(\cdots \otimes a \otimes c \otimes b \otimes \cdots)+\theta(\cdots \otimes a \otimes[b c] \otimes \cdots) \\
= & \theta(\cdots \otimes c \otimes a \otimes b \otimes \cdots)+\theta(\cdots \otimes[a c] \otimes b \otimes \cdots) \\
& +\theta(\cdots \otimes a \otimes[b c] \otimes \cdots) \\
= & \theta(\cdots \otimes c \otimes b \otimes a \otimes \cdots)+\theta(\cdots \otimes c \otimes[a b] \otimes \cdots) \\
& +\theta(\cdots \otimes b \otimes[a c] \otimes \cdots)+\theta(\cdots \otimes[[a c] b] \otimes \cdots) \\
& +\theta(\cdots \otimes a \otimes[b c] \otimes \cdots)
\end{aligned}
$$

again using the inductive assumptions. Comparing the two expressions obtained we see that they are equal since

$$
[[a c] b]=[[a b] c]+[[b c] a] .
$$

Thus $\theta: T^{n, i+1} \rightarrow R$ is now defined and this gives

$$
\theta: T^{0} \oplus \cdots \oplus T^{n-1} \oplus T^{n, i+1} \rightarrow R
$$

Since $T^{n}=T^{n, r}$ for $r$ sufficiently large we have

$$
\theta: T^{0} \oplus \cdots \oplus T^{n} \rightarrow R .
$$

Since $T=T^{0} \oplus T^{1} \oplus T^{2} \oplus \cdots$ we have defined $\theta: T \rightarrow R$ satisfying the required conditions.

We now return to part (b) of the proof of Theorem 9.4. We have $\mathfrak{H}(L)=T / J$ and the elements

$$
\begin{aligned}
x_{i_{1}} \otimes \cdots \otimes & x_{i_{k}} \otimes x_{i_{k+1}} \otimes \cdots \otimes x_{i_{n}}-x_{i_{1}} \otimes \cdots \otimes x_{i_{k+1}} \otimes x_{i_{k}} \otimes \cdots \otimes x_{i_{n}} \\
& -x_{i_{1}} \otimes \cdots \otimes\left[x_{i_{k}} x_{i_{k+1}}\right] \otimes \cdots \otimes x_{i_{n}} \\
= & x_{i_{1}} \ldots x_{i_{k-1}}\left(x_{i_{k}} \otimes x_{i_{k+1}}-x_{i_{k+1}} \otimes x_{i_{k}}-\left[x_{i_{k}} x_{i_{k+1}}\right]\right) x_{i_{k+2}} \ldots x_{i_{n}}
\end{aligned}
$$

all lie in $J$. In fact the definition of $J$ shows that each element of $J$ is a linear combination of such elements. Thus the linear map $\theta: T \rightarrow R$ of Lemma 9.5 annihilates all elements of $J$, and so induces a linear map $\bar{\theta}: T / J \rightarrow R$, that is $\bar{\theta}: \mathfrak{U}(L) \rightarrow R$. Now the monomial $y_{i_{1}}^{r_{1}} \ldots y_{i_{n}}^{r_{n}} \in \mathfrak{H}(L)$ for $i_{1}<\cdots<i_{n}$ is mapped by $\bar{\theta}$ to $z_{i_{1}}^{r_{1}} \ldots z_{i_{n}}^{r_{n}} \in R$. Since the elements $z_{i_{1}}^{r_{1}} \ldots z_{i_{n}}^{r_{n}}$ are linearly independent in the polynomial ring $R$ it follows that the elements $y_{i_{1}}^{r_{1}} \ldots y_{i_{n}}^{r_{n}}$ given in the statement of Theorem 9.4 must be linearly independent in $\mathfrak{l}(L)$. This completes the proof.

We now deduce some consequences of the Poincaré-Birkhoff-Witt basis theorem. (We shall subsequently call it the PBW basis theorem.)

Corollary 9.6 The map $\sigma: L \rightarrow \mathfrak{l}(L)$ is injective.
Proof. The elements $x_{i}, i \in I$, form a basis for $L$ and $\sigma\left(x_{i}\right)=y_{i}$. By the PBW basis theorem the elements $y_{i}, i \in I$, are linearly independent. Thus the kernel of $\sigma$ is zero.

Corollary 9.7 The subspace $\sigma(L)$ is a Lie subalgebra of $[\mathfrak{H}(L)]$ isomorphic to $L$. Thus $\sigma$ identifies $L$ with a Lie subalgebra of $[\mathfrak{H}(L)]$.

Proof. By Corollary 9.6 we know that $\sigma: L \rightarrow \sigma(L)$ is bijective. The elements $y_{i}, i \in I$, form a basis of $\sigma(L)$ and we have

$$
y_{i} y_{j}-y_{j} y_{i}=\sigma\left[x_{i} x_{j}\right] .
$$

It follows that $\left[y_{i} y_{j}\right] \in \sigma(L)$ and so $\sigma(L)$ is a Lie subalgebra of $[U(L)]$.

It is often convenient to consider $L$ as a subspace of $\mathfrak{H}(L)$ without mentioning the map $\sigma$ explicitly.

Corollary $9.8 \mathfrak{U}(L)$ has no zero-divisors.

Proof. Let $a, b \in \mathfrak{H}(L)$ have $a \neq 0, b \neq 0$. Then we have

$$
\begin{aligned}
a & =\sum \lambda_{i_{1}, \ldots, i_{n}, r_{1}, \ldots, r_{n}} y_{i_{1}}^{r_{1}} \ldots y_{i_{n}}^{r_{n}} \\
b & =\sum \mu_{i_{1}, \ldots, i_{n}, r_{1}, \ldots, r_{n}} y_{i_{1}}^{r_{1}} \ldots y_{i_{n}}^{r_{n}} .
\end{aligned}
$$

We write

$$
a=f\left(y_{i}\right)+\text { a sum of terms of smaller degree }
$$

where $f\left(y_{i}\right)$ is the sum of all terms $\lambda_{i_{1}, \ldots, i_{n}, r_{1}, \ldots, r_{n}} y_{i_{1}}^{r_{1}} \ldots y_{i_{n}}^{r_{n}}$ of maximal total degree $r=r_{1}+\cdots+r_{n}$. Similarly we have

$$
b=g\left(y_{i}\right)+\text { a sum of terms of smaller degree. }
$$

Now we have

$$
y_{i} y_{j}=y_{j} y_{i}+\text { a sum of terms of degree } 1
$$

and so

$$
f\left(y_{i}\right) g\left(y_{i}\right)=(f g)\left(y_{i}\right)+\text { a sum of terms of smaller degree. }
$$

Hence

$$
a b=(f g)\left(y_{i}\right)+\mathrm{a} \text { sum of terms of smaller degree. }
$$

Now $f$ is not the zero polynomial since $a \neq 0$ and $g$ is not the zero polynomial since $b \neq 0$. Thus $f g$ is not the zero polynomial. The PBW basis theorem then implies that $a b \neq 0$.

### 9.3 Free Lie algebras

It is well known how to define groups by generators and relations. One first constructs the free group on the given set of generators and then forms the factor group with respect to the smallest normal subgroup containing the elements specified by the given relations. We shall show that something similar can be done in the theory of Lie algebras. We first introduce the idea of the free Lie algebra $F L(X)$ on a set $X$.

Let $X=\left\{x_{i}, i \in I\right\}$ be a set of elements parametrised by an index set $I$. We first define the free associative algebra $F(X)$ on the set $X . F(X)$ is the set of all finite sums of the form

$$
\sum_{k \geq 0} \sum_{i_{1}, \ldots, i_{k} \in I} \lambda_{i_{1}, \ldots, i_{k}} x_{i_{1}} \ldots x_{i_{k}}
$$

with $\lambda_{i_{1}}, \ldots, i_{k} \in \mathbb{C}$, summed over all non-negative integers $k$ and all ordered $k$-tuples $i_{1}, \ldots, i_{k}$ from $I$ (repetitions being allowed). When $k=0$ the product $x_{i_{1}} \ldots x_{i_{k}}$ is the empty product, and is written as 1 . The operations of addition, multiplication and scalar multiplication are defined in an obvious way and make $F(X)$ into an associative algebra over $\mathbb{C}$ with identity 1.

Let $[F(X)]$ be the Lie algebra obtained from the associative algebra $F(X)$ in the usual manner. $X$ is a subset of $[F(X)]$. We define $F L(X)$ to be the intersection of all the Lie subalgebras of $[F(X)]$ containing $X$, i.e. the Lie subalgebra of $[F(X)]$ generated by $X . F L(X)$ is called the free Lie algebra on the set $X$. It is clear that $X$ is contained in $F L(X)$ so we have an injective $\operatorname{map} i: X \rightarrow F L(X)$.

In order to justify its name, we show that the free Lie algebra $F L(X)$ has the following universal property.

Proposition 9.9 Let $\theta: X \rightarrow L$ be any map from the set $X$ into a Lie algebra $L$. Then there is a unique homomorphism $\phi: F L(X) \rightarrow L$ such that $\phi \circ i=\theta$.


Proof. Consider the maps $X \xrightarrow{\theta} L \xrightarrow{\sigma} \mathfrak{H}(L)$. Let $\theta^{\prime}: X \rightarrow \mathfrak{H}(L)$ be given by $\theta^{\prime}=\sigma \circ \theta$. The map $\theta^{\prime}$ from $X$ into $\mathfrak{H}(L)$ can be extended uniquely (in an obvious way) to an associative algebra homomorphism $\phi^{\prime}: F(X) \rightarrow \mathfrak{U}(L)$. The same map gives a Lie algebra homomorphism $\phi^{\prime}:[F(X)] \rightarrow[\mathfrak{U}(L)]$. Now we have $\phi^{\prime}(X) \subset \sigma(L)$ and we know from Corollary 9.7 that $\sigma(L)$ is a Lie subalgebra of $[\mathfrak{U}(L)]$ isomorphic to $L$. The set of elements of $[F(X)]$ mapped by $\phi^{\prime}$ into $\sigma(L)$ is therefore a Lie subalgebra of $[F(X)]$ containing $X$, and this contains $F L(X)$. Hence we have $\phi^{\prime}: F L(X) \rightarrow \sigma(L)$. We define $\phi: F L(X) \rightarrow L$ by $\phi=\sigma^{-1} \circ \phi^{\prime}$. We check that $\phi \circ i=\theta$. For if $x \in X$ we have

$$
\phi \circ i(x)=\sigma^{-1} \phi^{\prime} i(x)=\sigma^{-1} \phi^{\prime}(x)=\sigma^{-1} \theta^{\prime}(x)=\theta(x) .
$$

Thus we have a homomorphism $\phi$ of the required type. Finally we show that $\phi$ is unique. Let $\bar{\phi}: F L(X) \rightarrow L$ be another such homomorphism. Then we have

$$
\phi i(x)=\theta(x)=\bar{\phi} i(x) \quad \text { for all } x \in X
$$

Thus $\phi$ agrees with $\bar{\phi}$ on $X$. Now the set of elements of $F L(X)$ for which $\phi$ agrees with $\bar{\phi}$ is a Lie subalgebra of $F L(X)$ containing $X$. Since $X$ generates $F L(X)$ as a Lie algebra we deduce that $\phi$ agrees with $\bar{\phi}$ on $F L(X)$.

We next identify the universal enveloping algebra of the free Lie algebra $F L(X)$. This turns out to be isomorphic to the free associative algebra $F(X)$.

Proposition 9.10 The universal enveloping algebra $\mathfrak{H}(F L(X))$ is isomorphic to $F(X)$.

Proof. We have an inclusion map $\sigma: F L(X) \rightarrow F(X)$. We shall show that the universal property of enveloping algebras given in Proposition 9.2 is satisfied by $F(X)$. Thus we shall show that if $A$ is any associative algebra with 1 over $\mathbb{C}$ and if $\theta: F L(X) \rightarrow[A]$ is any Lie algebra homomorphism then there exists a unique associative algebra homomorphism $\phi: F(X) \rightarrow A$ such that $\phi \circ \sigma=\theta$.

Now the Lie homomorphism $\theta: F L(X) \rightarrow[A]$ restricts to a map $\theta: X \rightarrow A$. This map can be extended to a unique associative algebra homomorphism $\phi: F(X) \rightarrow A$. This same map gives a Lie algebra homomorphism $\phi:[F(X)] \rightarrow[A]$. By restriction we obtain a Lie algebra homomorphism $\phi: F L(X) \rightarrow[A]$. However, $\phi$ agrees with $\theta$ on $X$ and $X$ generates $F L(X)$ as a Lie algebra. Hence $\phi$ agrees with $\theta$ on $F L(X)$. It follows that $\phi \circ \sigma=\theta$ as required. Thus there exists an algebra homomorphism $\phi$ of the required kind. On the other hand $\phi$ is clearly unique since $F L(X)$ contains $X$ and therefore generates the associative algebra $F(X)$.

Thus $F(X)$ satisfies the above universal property. Of course $\mathfrak{H}(F L(X))$ satisfies it also. This implies that $\mathfrak{H}(F L(X))$ is isomorphic to $F(X)$. For suppose we are given a Lie algebra $L$ and two associative algebras $\mathfrak{U}, \mathfrak{U}^{\prime}$ with maps $\sigma: L \rightarrow \mathfrak{U}, \sigma^{\prime}: L \rightarrow \mathfrak{U}^{\prime}$ both satisfying the universal property. Then we obtain unique algebra homomorphisms $\phi: \mathfrak{U} \rightarrow \mathfrak{U}^{\prime}$ and $\phi^{\prime}: \mathfrak{U}^{\prime} \rightarrow \mathfrak{U}$ such that $\sigma^{\prime}=\phi \circ \sigma$ and $\sigma=\phi^{\prime} \circ \sigma^{\prime}$.


It follows that

$$
\phi^{\prime} \phi \sigma(x)=\sigma(x), \quad \phi \phi^{\prime} \sigma^{\prime}(x)=\sigma^{\prime}(x)
$$

for all $x \in L$. Now $\sigma(L)$ generates $\mathfrak{U}$ and $\sigma^{\prime}(L)$ generates $\mathfrak{U}^{\prime}$ as associative algebras, by the uniqueness of $\phi$ and $\phi^{\prime}$. It follows that

$$
\phi^{\prime} \phi=\mathrm{Id}_{\mathfrak{l}}, \quad \phi \phi^{\prime}=\mathrm{Id}_{\mathfrak{l}^{\prime}}
$$

and so $\phi, \phi^{\prime}$ are inverse isomorphisms between $\mathfrak{U}$ and $\mathfrak{U}^{\prime}$.

### 9.4 Lie algebras defined by generators and relations

Let $X=\left\{x_{i}, i \in I\right\}$ be a given set. A Lie monomial in the elements of $X$ is a finite product of elements of $X$ bracketed by Lie brackets in any manner. For example

$$
\left[\left[\left[x_{3}\left[x_{1} x_{2}\right]\right] x_{3}\right]\left[x_{2}\left[x_{1} x_{1}\right]\right]\right]
$$

is a Lie monomial on the set $X=\left\{x_{1}, x_{2}, x_{3}\right\}$. A Lie word in the elements on $X$ is a finite linear combination of Lie monomials on $X$ with coefficients in $\mathbb{C}$. For example

$$
3\left[\left[\left[x_{3}\left[x_{1} x_{2}\right]\right] x_{3}\right]\left[x_{2}\left[x_{1} x_{1}\right]\right]\right]+2\left[x_{2}\left[\left[x_{1} x_{2}\right]\left[x_{3} x_{2}\right]\right]\right]
$$

is a Lie word on the set $X=\left\{x_{1}, x_{2}, x_{3}\right\}$.
Let $R=\left\{w_{j}, j \in J\right\}$ be a set of Lie words in the elements of $X$. We shall define a Lie algebra $L(X ; R)$ called the Lie algebra generated by $X$ subject to relations $R$.

Now the elements of $X$ all lie in the free Lie algebra $F L(X)$ and all the Lie words $w_{j}$ also lie in $F L(X)$. Of course different Lie words can give the same element of $F L(X)$ because of relations such as $\left[x_{i} x_{i}\right]=0$ and the Jacobi identity. Let $\langle R\rangle$ be the ideal of $F L(X)$ generated by $R$. Thus $\langle R\rangle$ is the intersection of all ideals of $F L(X)$ containing $R$. We define $L(X ; R)$ by

$$
L(X ; R)=F L(X) /\langle R\rangle
$$

Lemma 9.11 Let $R, R^{\prime}$ be sets of Lie words in $X$ such that $R^{\prime} \subset R$. Then $L(X ; R)$ is isomorphic to a factor algebra of $L\left(X ; R^{\prime}\right)$.

Proof. Since $R^{\prime} \subset R$ we have

$$
\left\langle R^{\prime}\right\rangle \subset\langle R\rangle \subset F L(X)
$$

It follows that

$$
L(X ; R)=\frac{F L(X)}{\langle R\rangle} \cong \frac{F L(X)}{\left\langle R^{\prime}\right\rangle} / \frac{\langle R\rangle}{\left\langle R^{\prime}\right\rangle}=\frac{L\left(X ; R^{\prime}\right)}{I}
$$

where $I=\langle R\rangle /\left\langle R^{\prime}\right\rangle$.
Example 9.12 Let $A$ be a Cartan matrix on the standard list 6.12. In Section 7.4 we defined a Lie algebra $L(A)$ associated with $A$, and $L(A)$ was subsequently shown in Proposition 7.35 to be a simple Lie algebra. In fact all the finite dimensional non-trivial simple Lie algebras over $\mathbb{C}$ have form $L(A)$, as $A$ runs over all Cartan matrices on the standard list. The definition of $L(A)$ given in Section 7.4 shows that $L(A)$ can conveniently be described in terms of generators and relations. In fact we have

$$
L(A) \cong L(X ; R)
$$

where $X=\left\{e_{1}, \ldots, e_{l}, h_{1}, \ldots, h_{l}, f_{1}, \ldots, f_{l}\right\}$ and $R$ is the set of Lie words in $X$ given by

$$
\begin{array}{ll}
{\left[h_{i} h_{j}\right]} & \\
{\left[h_{i} e_{j}\right]-A_{i j} e_{j}} & \\
{\left[h_{i} f_{j}\right]+A_{i j} f_{j}} & \\
{\left[e_{i} f_{i}\right]-h_{i}} & \\
{\left[e_{i} f_{j}\right] \quad \text { for } i \neq j} & \\
{\left[e_{i}\left[e_{i}\left[\ldots\left[e_{i} e_{j}\right]\right]\right]\right]} & \text { for } i \neq j \\
{\left[f_{i}\left[f_{i}\left[\ldots\left[f_{i} f_{j}\right]\right]\right]\right]} & \text { for } i \neq j
\end{array}
$$

where the number of occurrences of $e_{i}, f_{i}$ respectively in the last two words is $1-A_{i j}$.

Example 9.13 Again let $A$ be a Cartan matrix on the standard list 6.12. In Section 7.4 we also defined a certain Lie algebra $\tilde{L}(A)$ depending on $A$ which contains $L(A)$ as a factor algebra. The algebra $\tilde{L}(A)$ is infinite dimensional. It can also conveniently be described by generators and relations. In fact we have

$$
\tilde{L}(A) \cong L\left(X ; R^{\prime}\right)
$$

where $X=\left\{e_{1}, \ldots, e_{l}, h_{1}, \ldots, h_{l}, f_{1}, \ldots, f_{l}\right\}$ and $R^{\prime}$ is the set of Lie words on $X$ given by

$$
\begin{aligned}
& {\left[h_{i} h_{j}\right]} \\
& {\left[h_{i} e_{j}\right]-A_{i j} e_{j}} \\
& {\left[h_{i} f_{j}\right]+A_{i j} f_{j}} \\
& {\left[e_{i} f_{i}\right]-h_{i}} \\
& {\left[e_{i} f_{j}\right] \quad \text { for } i \neq j}
\end{aligned}
$$

We observe that $R^{\prime}$ is a proper subset of the set $R$ of relations in Example 9.12. This explains why $L(A)$ is isomorphic to a factor algebra of $\tilde{L}(A)$, as in Lemma 9.11.

### 9.5 Graph automorphisms of simple Lie algebras

Let $A$ be a Cartan matrix on the standard list 6.12 and $\sigma$ be a permutation of $\{1, \ldots, l\}$ such that $A_{\sigma(i) \sigma(j)}=A_{i j}$ for all $i, j$. Let $L(A)$ be the simple Lie algebra associated with $A . L(A)$ can be generated by $e_{1}, \ldots, e_{l}, h_{1}, \ldots, h_{l}, f_{1}, \ldots, f_{l}$. We define a permutation of this generating set by

$$
e_{i} \rightarrow e_{\sigma(i)} \quad f_{i} \rightarrow f_{\sigma(i)} \quad h_{i} \rightarrow h_{\sigma(i)}
$$

Under this permutation of the generators each of the defining relations of $L(A)$ in Example 9.12 is transformed into a defining relation. Let

$$
L(A)=L(X ; R)=\frac{F L(X)}{\langle R\rangle}
$$

The given permutation of $X$ extends to a Lie algebra homomorphism of $F L(X)$ into itself, and this homomorphism maps the ideal $\langle R\rangle$ into itself. It therefore induces a Lie algebra homomorphism of $L(X ; R)$ into itself. Since the permutation of $X$ is invertible, so is this Lie algebra homomorphism. It is thus an isomorphism of $L(X ; R)$ into itself, that is an automorphism of $L(A)$. This automorphism is called a graph automorphism of $L(A)$ and will also be denoted by $\sigma$. The possible non-trivial graph automorphisms can be described in terms of the action of $\sigma$ on the Dynkin diagram of $L(A)$. These possibilities are listed below.

Type $A_{2 k}$



Type $D_{k+1}$


$$
\begin{aligned}
& \sigma(i)=i \text { for } 1 \leq i \leq k-1 \\
& \sigma(k)=k+1 \\
& \sigma(k+1)=k
\end{aligned}
$$

Type $D_{4}$


$$
\sigma(1)=1 \quad \sigma(2)=3 \quad \sigma(3)=4 \quad \sigma(4)=2
$$

(The inverse of $\sigma$ is also a graph automorphism, which can be obtained from $\sigma$ by renumbering the vertices.)

Type $E_{6}$


$$
\sigma(1)=6 \quad \sigma(2)=5 \quad \sigma(3)=3 \quad \sigma(4)=4 \quad \sigma(5)=2 \quad \sigma(6)=1
$$

Our main aim in the present section is to determine the fixed point subalgebra

$$
L(A)^{\sigma}=\{x \in L(A) ; \sigma(x)=x\} .
$$

We begin by considering the action of $\sigma$ on $V=H_{\mathbb{R}}^{*}$ given by $\sigma\left(\alpha_{i}\right)=\alpha_{\sigma(i)}$ and extending by linearity. Let $V^{1}=\{v \in V ; \sigma(v)=v\}$. For each orbit $J$ of $\sigma$ on $\{1, \ldots, l\}$ we define $\alpha_{J}=\frac{1}{|J|} \sum_{j \in J} \alpha_{j}$. Then $\alpha_{J} \in V^{1}$ and the $\alpha_{J}$ form a basis of $V^{1}$ as $J$ runs over the $\sigma$-orbits on $\{1, \ldots, l\} . \alpha_{J}$ is simply the projection of $\alpha_{j}$ on to the subspace $V^{1}$ of the Euclidean space $V$. We see from the above diagrams that the orbits $J$ have the following possible types.
(a) $|J|=1 \quad$ and $\quad J=\{j\} \quad$ with $\sigma(j)=j$.
(b) $|J|=2 \quad$ and $\quad J=\{j \bar{j}\} \quad$ where $\sigma(j)=\bar{j}, \sigma(\bar{j})=j \quad$ and $\quad \alpha_{j}+\alpha_{\bar{j}} \notin \Phi$.
(c) $|J|=3 \quad$ and $\quad J=\{j \bar{j} \overline{\bar{j}}\} \quad$ where $\quad \sigma(j)=\bar{j}, \sigma(\bar{j})=\overline{\bar{j}}, \sigma(\overline{\bar{j}})=j \quad$ and $\alpha_{j}+\alpha_{\bar{j}}, \alpha_{j}+\alpha_{\bar{j}}, \alpha_{\bar{j}}+\alpha_{\bar{j}}$ do not lie in $\Phi$.
(d) $|J|=2 \quad$ and $\quad J=\{j \bar{j}\} \quad$ where $\sigma(j)=\bar{j}, \sigma(\bar{j})=j \quad$ and $\quad \alpha_{j}+\alpha_{\bar{j}} \in \Phi$.

These four will be called orbits of types $A_{1}, A_{1} \times A_{1}, A_{1} \times A_{1} \times A_{1}$ and $A_{2}$ respectively.

We next consider the possible pairs $J, K$ of distinct orbits.

Lemma 9.14 The vectors $\alpha_{J}, \alpha_{K}$ for distinct $\sigma$-orbits $J$, $K$ form a fundamental system of roots of rank 2. The type of this rank 2 system is as follows.

$$
\mathbf{J} \quad \mathbf{K} \quad \text { Type of }\left\{\alpha_{J}, \alpha_{k}\right\}
$$

(i)

(ii)

(iii)

(iv)

(v)

(vi) If no node in $J$ is joined to any node in $K$ then the type of $\left\{\alpha_{j}, \alpha_{k}\right\}$ is $A_{1} \times A_{1}$.

Proof. This is straightforward. Suppose for example we have case (v) with roots numbered


Then

$$
\alpha_{J}=\frac{\alpha_{1}+\alpha_{4}}{2}, \quad \alpha_{K}=\frac{\alpha_{2}+\alpha_{3}}{2}
$$

We have

$$
\begin{aligned}
\left\langle\alpha_{J}, \alpha_{J}\right\rangle & =\frac{1}{2}\left\langle\alpha_{1}, \alpha_{1}\right\rangle \\
\left\langle\alpha_{K}, \alpha_{K}\right\rangle & =\frac{1}{4}\left\langle\alpha_{1}, \alpha_{1}\right\rangle \\
\left\langle\alpha_{J}, \alpha_{K}\right\rangle & =-\frac{1}{4}\left\langle\alpha_{1}, \alpha_{1}\right\rangle .
\end{aligned}
$$

Thus $\left\langle\alpha_{J}, \alpha_{J}\right\rangle=2\left\langle\alpha_{K}, \alpha_{K}\right\rangle$ and $2\left\langle\alpha_{J}, \alpha_{K}\right\rangle /\left\langle\alpha_{J}, \alpha_{J}\right\rangle=-1$.
Hence we have a fundamental system with diagram


Corollary 9.15 Let $\Pi^{1}$ be the set of vectors $\alpha_{J}$ for all $\sigma$-orbits J on $\{1, \ldots, l\}$. Then $\Pi^{1}$ is a fundamental system of roots of the following type:

| Type $\Pi$ | Order of $\sigma$ | Type $\Pi^{1}$ |
| :--- | :---: | :---: |
| $A_{2 k}$ | 2 | $B_{k}$ |
| $A_{2 k-1}$ | 2 | $C_{k}$ |
| $D_{k+1}$ | 2 | $B_{k}$ |
| $D_{4}$ | 3 | $G_{2}$ |
| $E_{6}$ | 2 | $F_{4}$ |

Proof. This follows immediately from Lemma 9.14.
The relationship between $\Pi$ and $\Pi^{1}$ may be illustrated in the following diagrams.

п

$\Pi^{1}$

$\Pi$

$\Pi^{1}$
 -•

$\Pi$


- . .

$\Pi^{1}$

$\Pi$


Now let $\Phi^{1}$ be the root system in $V^{1}$ with fundamental system $\Pi^{1}$. Let $W^{1}$ be the Weyl group of $\Phi^{1}$. Then $\Phi^{1}=W^{1}\left(\Pi^{1}\right)$. Let $A^{1}$ be the Cartan matrix of $\Phi^{1}$.

Proposition 9.16 Let $I, J$ be distinct $\sigma$-orbits on $\{1, \ldots, l\}$. Then
$A_{I J}^{1}= \begin{cases}\sum_{i \in I} A_{i j} & \text { for any } j \in J, \text { if I has type } A_{1}, A_{1} \times A_{1} \text { or } A_{1} \times A_{1} \times A_{1} \\ 2 \sum_{i \in I} A_{i j} & \text { for any } j \in J, \text { if I has type } A_{2} .\end{cases}$

Proof. This follows from Lemma 9.14.

Proposition 9.17 Let $W^{\sigma}=\{w \in W ; w \sigma=\sigma w$ on $V\}$. Then there is an isomorphism $W^{1} \rightarrow W^{\sigma}$ under which the fundamental reflection $s_{J} \in W^{1}$
corresponding to $\alpha_{J}$ maps to $\left(w_{0}\right)_{J} \in W^{\sigma}$, the element of maximal length in the Weyl subgroup $W_{J}$ of $W$ generated by the $s_{i}$ for $i \in J$.

Proof. We first observe that $W^{\sigma}$ acts on $V^{1}$. For let $w \in W^{\sigma}, v \in V^{1}$. Then

$$
\sigma(w v)=w(\sigma v)=w v, \quad \text { thus } w v \in V^{1} .
$$

Secondly we note that $\left(w_{0}\right)_{J} \in W^{\sigma}$ for each $\sigma$-orbit $J$. For $j \in J$ we have

$$
\sigma s_{j} \sigma^{-1}=s_{\sigma(j)}
$$

thus $\sigma W_{J} \sigma^{-1}=W_{J}$. Since $\sigma$ preserves the sign of each root we have

$$
\sigma\left(w_{0}\right)_{J} \sigma^{-1}\left(\Phi_{J}^{+}\right)=\Phi_{J}^{-}
$$

and hence

$$
\sigma\left(w_{0}\right)_{J} \sigma^{-1}=\left(w_{0}\right)_{J}
$$

by Proposition 5.17. Thus $\left(w_{0}\right)_{J} \in W^{\sigma}$.
Thirdly we note that the element $\left(w_{0}\right)_{J} \in W^{\sigma}$, when restricted to $V^{1}$, coincides with $s_{J}$. For

$$
\left(w_{0}\right)_{J}\left(\alpha_{J}\right)=\left(w_{0}\right)_{J}\left(\frac{1}{|J|} \sum_{j \in J} \alpha_{j}\right)=-\frac{1}{|J|} \sum_{j \in J} \alpha_{j}=-\alpha_{J}
$$

since $\left(w_{0}\right)_{J}\left(\Phi_{J}^{+}\right)=\Phi_{J}^{-}$.
Also if $v \in V^{1}$ satisfies $\left\langle\alpha_{J}, v\right\rangle=0$ then it satisfies $\left\langle\alpha_{j}, v\right\rangle=0$ for all $j \in J$. It follows that $\left(w_{0}\right)_{J}(v)=v$. Thus $\left(w_{0}\right)_{J}$ coincides with $s_{J}$ on restriction to $V^{1}$.

We next show that the elements $\left(w_{0}\right)_{J}$ for all $\sigma$-orbits $J$ generate $W^{\sigma}$. Let $w \in W^{\sigma}$ satisfy $w \neq 1$. Then there exists a fundamental root $\alpha_{j}$ with $w\left(\alpha_{j}\right) \in \Phi^{-}$. Let $J$ be the $\sigma$-orbit containing $j$. Then

$$
w \sigma\left(\alpha_{j}\right)=\sigma w\left(\alpha_{j}\right) \in \Phi^{-}
$$

since $\sigma$ preserves the sign of each root. Thus $w\left(\alpha_{i}\right) \in \Phi^{-}$for all $i \in J$. Now $\left(w_{0}\right)_{J}$ changes the signs of all roots in $\Phi_{J}$ but of none in $\Phi-\Phi_{J}$. Hence

$$
l\left(w\left(w_{0}\right)_{J}\right)=l(w)-l\left(\left(w_{0}\right)_{J}\right)<l(w) .
$$

We assume by induction on $l(w)$ that $w\left(w_{0}\right)_{J}$ lies in the subgroup generated by the $\left(w_{0}\right)_{I}$ for all $\sigma$-orbits $I$. It follows that $w$ has the same property. Hence the $\left(w_{0}\right)_{I}$ generate $W^{\sigma}$.

We may now define a homomorphism $W^{\sigma} \rightarrow W^{1}$, by restricting the action of $w \in W^{\sigma}$ from $V$ to $V^{1}$. Since $W^{\sigma}$ is generated by the elements $\left(w_{0}\right)_{J}$ and $\left(w_{0}\right)_{J}$ restricted to $V^{1}$ is $s_{J}$, the image of the homomorphism is generated by
the $s_{J}$ and so is $W^{1}$. Finally we show our map is injective. Suppose $w \in W^{\sigma}$ and $w \neq 1$. Then there exists a $\sigma$-orbit $J$ such that $w\left(\alpha_{i}\right) \in \Phi^{-}$for all $i \in J$. Thus $w\left(\alpha_{J}\right) \neq \alpha_{J}$ and so $w$ acts non-trivially on $V^{1}$. Thus our map $W^{\sigma} \rightarrow W^{1}$ is an isomorphism under which $\left(w_{0}\right)_{J} \in W^{\sigma}$ corresponds to $s_{J} \in W^{1}$.

We next consider the relation between the root systems $\Phi$ and $\Phi^{1}$. For each $\alpha \in \Phi$ we denote by $\alpha^{1}$ its projection into $V^{1}$.

Proposition 9.18 (a) For each $\alpha \in \Phi, \alpha^{1}$ is a positive multiple of a root in $\Phi^{1}$.
(b) Let $\sim$ be the equivalence relation on $\Phi$ given by $\alpha \sim \beta$ if and only if $\alpha^{1}$ is a positive multiple of $\beta^{1}$. Then the equivalence classes are the subsets of $\Phi$ of form $w\left(\Phi_{J}^{+}\right)$where $J$ is a $\sigma$-orbit in $\{1, \ldots, l\}$ and $w \in W^{\sigma}$.
(c) There is a bijection between equivalence classes on $\Phi$ and roots in $\Phi^{1}$ given by $w\left(\Phi_{J}^{+}\right) \leftrightarrow w^{1}\left(\alpha_{J}\right)$ where $w^{1}$ is the restriction of $w$ to $V^{1}$.

Proof. We first show that each $\alpha \in \Phi$ lies in $w\left(\Phi_{J}^{+}\right)$for some $\sigma$-orbit $J$ and some $w \in W^{\sigma}$. Consider the element $w_{0} \in W$ of maximal length. $w_{0}$ transforms each positive root to a negative root. Since $\sigma$ does not change the sign of any root we have

$$
\sigma w_{0} \sigma^{-1}\left(\Phi^{+}\right)=\Phi^{-}
$$

Since $\sigma w_{0} \sigma^{-1} \in W$ it follows that $\sigma w_{0} \sigma^{-1}=w_{0}$, that is $w_{0} \in W^{1}$. By Proposition 9.17 the elements $\left(w_{0}\right)_{J}$ for all $\sigma$-orbits $J$ generate $W^{\sigma}$ and so

$$
w_{0}=\left(w_{0}\right)_{J_{1}} \ldots\left(w_{0}\right)_{J_{r}}
$$

for some $J_{1}, \ldots, J_{r}$. Let $\alpha \in \Phi^{+}$. Then $w_{0}(\alpha) \in \Phi^{-}$. Thus there exists $i$ such that

$$
\begin{aligned}
& \left(w_{0}\right)_{J_{i+1}} \ldots\left(w_{0}\right)_{J_{r}}(\alpha) \in \Phi^{+} \\
& \left(w_{0}\right)_{J_{i}}\left(w_{0}\right)_{J_{i+1}} \ldots\left(w_{0}\right)_{J_{r}}(\alpha) \in \Phi^{-}
\end{aligned}
$$

Since the only positive roots made negative by $\left(w_{0}\right)_{J_{i}}$ are those in $\Phi_{J_{i}}^{+}$we have

$$
\left(w_{0}\right)_{J_{i+1}} \ldots\left(w_{0}\right)_{J_{r}}(\alpha) \in \Phi_{J_{i}}^{+}
$$

that is $\alpha \in\left(w_{0}\right)_{J_{r}} \ldots\left(w_{0}\right)_{J_{i+1}}\left(\Phi_{J_{i}}^{+}\right)$and $-\alpha \in\left(w_{0}\right)_{J_{r}} \ldots\left(w_{0}\right)_{J_{l+1}}\left(w_{0}\right)_{J_{i}}\left(\Phi_{J_{i}}^{+}\right)$. Hence each root in $\Phi$ lies in $w\left(\Phi_{J}^{+}\right)$for some $\sigma$-orbit $J$ and some $w \in W^{\sigma}$.

Now consider the projection $\alpha^{1}$ for $\alpha \in \Phi_{J}^{+}$. If $J$ has type $A_{1}, A_{1} \times A_{1}$ or $A_{1} \times A_{1} \times A_{1}$ then $\Phi_{J}^{+}=\Pi_{J}$ and so $\alpha^{1}=\alpha_{J}$ for $\alpha \in \Phi_{J}^{+}$. If $J$ has type $A_{2}$, however, then $\Pi_{J}=\left\{\alpha_{j}, \alpha_{\bar{j}}\right\}$ and $\Phi_{J}^{+}=\left\{\alpha_{j}, \alpha_{\bar{j}}, \alpha_{j}+\alpha_{\bar{j}}\right\}$. We have

$$
\alpha^{1}= \begin{cases}\alpha_{J} & \text { when } \alpha=\alpha_{j} \text { or } \alpha_{\bar{j}} \\ 2 \alpha_{J} & \text { when } \alpha=\alpha_{j}+\alpha_{\bar{j}}^{-}\end{cases}
$$

Thus $\alpha^{1}$ is a positive multiple of $\alpha_{J}$ when $\alpha \in \Phi_{J}^{+}$. Hence for $\alpha \in w\left(\Phi_{J}^{+}\right)$with $w \in W^{\sigma}$ we see that $\alpha^{1}$ is a positive multiple of $w\left(\alpha_{J}\right) \in \Phi^{1}$.

Now consider the equivalence relation on $\Phi$ defined in (b). The elements of each set $w\left(\Phi_{J}^{+}\right)$for $w \in W^{\sigma}$ lie in an equivalence class. Suppose $w\left(\Phi_{J}^{+}\right), w^{\prime}\left(\Phi_{K}^{+}\right)$lie in the same equivalence class for $\sigma$-orbits $J, K$ and $w, w \in W^{\sigma}$. Then

$$
w\left(\alpha_{J}\right)=w^{\prime}\left(\alpha_{K}\right) \in \Phi^{1}
$$

Hence $w^{\prime-1} w\left(\alpha_{J}\right)=\alpha_{K}$.
Consider the root $w^{\prime-1} w\left(\alpha_{j}\right) \in \Phi$ for $j \in J$. This root has the property that

$$
\left(w^{\prime-1} w\left(\alpha_{j}\right)\right)^{1}=\alpha_{K}
$$

Since $K$ is a $\sigma$-orbit this implies that $w^{\prime-1} w\left(\alpha_{j}\right)$ is a non-negative combination of the $\alpha_{k}$ for $k \in K$. Hence

$$
w^{\prime-1} w\left(\Pi_{J}\right) \subset \Phi_{K}^{+}
$$

and so $w^{\prime-1} w\left(\Phi_{J}^{+}\right) \subset \Phi_{K}^{+}$. By symmetry we also have

$$
w^{\prime} w^{-1}\left(\Phi_{K}^{+}\right) \subset \Phi_{J}^{+}
$$

Hence we have equality, that is

$$
w\left(\Phi_{J}^{+}\right)=w^{\prime}\left(\Phi_{K}^{+}\right)
$$

Hence the equivalence classes are the subsets of $\Phi$ of form $w\left(\Phi_{J}^{+}\right)$.
Now any root in $\Phi^{1}$ has form $w\left(\alpha_{J}\right)$ for some $w \in W^{1}$ and some $\sigma$-orbit $J$. The set of roots $\alpha \in \Phi$ such that $\alpha^{1}$ is a positive multiple of $w\left(\alpha_{J}\right)$ is $w\left(\Phi_{J}^{+}\right)$, as shown above. Thus

$$
w\left(\Phi_{J}^{+}\right) \leftrightarrow w\left(\alpha_{J}\right)
$$

is a bijective correspondence between equivalence classes on $\Phi$ and elements of $\Phi^{1}$.

Theorem 9.19 Let $\sigma$ be a graph automorphism of the simple Lie algebra $L(A)$. Then the subalgebra $L(A)^{\sigma}$ is isomorphic to the simple Lie algebra $L\left(A^{1}\right)$.

Proof. For each $\sigma$-orbit $J$ on $\{1, \ldots, l\}$ we define elements $e_{J}, h_{J}, f_{J}$ of $L(A)^{\sigma}$ by

$$
e_{J}=\sum_{j \in J} e_{j} \quad f_{J}=\sum_{j \in J} f_{j} \quad h_{J}=\sum_{j \in J} h_{j}
$$

if $J$ has type $A_{1}, A_{1} \times A_{1}$ or $A_{1} \times A_{1} \times A_{1}$ and

$$
e_{J}=\sqrt{ } 2 \sum_{j \in J} e_{j}, \quad f_{J}=\sqrt{ } 2 \sum_{j \in J} f_{j}, \quad h_{J}=2 \sum_{j \in J} h_{j}
$$

if $J$ has type $A_{2}$. Then we have

$$
\begin{aligned}
{\left[e_{I} f_{I}\right] } & =h_{I} \\
{\left[e_{I} f_{J}\right] } & =0 \quad \text { if } I \neq J \\
{\left[h_{I} h_{J}\right] } & =0
\end{aligned}
$$

We consider $\left[h_{I} e_{J}\right]$. If $I, J$ have type $A_{1}, A_{1} \times A_{1}$ or $A_{1} \times A_{1} \times A_{1}$ we have

$$
\left[h_{I} e_{J}\right]=\left[\sum_{i \in I} h_{i}, \sum_{j \in J} e_{j}\right]=\sum_{j}\left(\sum_{i} A_{i j} e_{j}\right)=A_{I J}^{1} e_{J}
$$

We also have $\left[h_{I} e_{J}\right]=A_{I J}^{1} e_{J}$ if one or both of $I, J$ has type $A_{2}$. Similarly

$$
\left[h_{I} f_{J}\right]=-A_{I J}^{1} f_{J} \quad \text { for all } I, J .
$$

We also check the relation

$$
\left[e_{I}\left[e_{I}\left[\ldots\left[e_{I} e_{J}\right]\right]\right]\right]=0 \quad \text { for } I \neq J
$$

where there are $1-A_{I J}^{1}$ factors $e_{I}$. This follows from the following observations, which can be checked from Lemma 9.14.

$$
\begin{aligned}
& \text { If } A_{I J}^{1}=0 \text { then }\left[e_{i} e_{j}\right]=0 \text { for all } i \in I, j \in J \\
& \text { If } A_{I J}^{1}=-1 \text { then }\left[e_{i}\left[e_{i^{\prime}} e_{j}\right]\right]=0 \text { for all } i, i^{\prime} \in I, j \in J \\
& \text { If } A_{I J}^{1}=-2 \text { then }\left[e_{i}\left[e_{i^{\prime}}\left[e_{i^{\prime \prime}} e_{j}\right]\right]\right]=0 \text { for all } i, i^{\prime}, i^{\prime \prime} \in I, j \in J . \\
& \text { If } A_{I J}^{1}=-3 \text { then }\left[e_{i}\left[e_{i^{\prime}}\left[e_{i^{\prime \prime}}\left[e_{i^{\prime \prime \prime}} e_{j}\right]\right]\right]\right]=0 \text { for all } i, i^{\prime}, i^{\prime \prime}, i^{\prime \prime \prime} \in I, j \in J .
\end{aligned}
$$

Similarly we obtain the relation

$$
\left[f_{I}\left[f_{I}\left[\cdots\left[f_{I} f_{J}\right]\right]\right]\right]=0 \quad \text { for } I \neq J
$$

with $1-A_{I J}^{1}$ factors $f_{I}$.

We now consider the generators and defining relations for the simple Lie algebra $L\left(A^{1}\right)$ given in Example 9.12. All these relations are satisfied by the elements $e_{J}, f_{J}, h_{J}$ of $L(A)^{\sigma}$. Thus there is a homomorphism

$$
L\left(A^{1}\right) \rightarrow L(A)^{\sigma}
$$

under which the generators of $L\left(A^{1}\right)$ map to the elements $e_{J}, f_{J}, h_{J}$ of $L(A)^{\sigma}$. Since $L\left(A^{1}\right)$ is simple this homomorphism is injective. We show it is also surjective and that the map is therefore an isomorphism. It will be sufficient to show that

$$
\operatorname{dim} L(A)^{\sigma}=\operatorname{dim} L\left(A^{1}\right)
$$

We consider the decomposition of $\Phi$ into equivalence classes given in Proposition 9.18. For each equivalence class $S$ let

$$
L_{S}=\bigoplus_{\alpha \in S} L_{\alpha}
$$

Then $\sigma\left(L_{S}\right)=L_{S}$ and

$$
\begin{aligned}
L & =H \oplus \sum_{S} L_{S} \\
L^{\sigma} & =H^{\sigma} \oplus \sum_{S}\left(L_{S}\right)^{\sigma}
\end{aligned}
$$

Now $\operatorname{dim}\left(L_{S}\right)^{\sigma} \leq 1$ for each equivalence class $S$. This is clear if $S$ has type $A_{1}, A_{1} \times A_{1}$ or $A_{1} \times A_{1} \times A_{1}$. Suppose then that $S$ has type $A_{2}$. Then $S=$ $\{\alpha, \beta, \alpha+\beta\}$. We have

$$
\sigma\left(e_{\alpha}\right)=\lambda e_{\beta}, \quad \sigma\left(e_{\beta}\right)=\lambda^{-1} e_{\alpha}
$$

for some $\lambda \in \mathbb{C}$. Hence

$$
\sigma\left[e_{\alpha} e_{\beta}\right]=\left[e_{\beta} e_{\alpha}\right]=-\left[e_{\alpha} e_{\beta}\right] .
$$

Thus $\left(L_{S}\right)^{\sigma}$ consists of all multiples of $e_{\alpha}+\lambda e_{\beta}$ and $\operatorname{dim}\left(L_{S}\right)^{\sigma}=1$. It follows that

$$
\begin{aligned}
\operatorname{dim} L^{\sigma} & \leq \operatorname{dim} H^{\sigma}+\text { no. of equivalence classes } S \\
& =\operatorname{dim} H^{1}+\left|\Phi^{1}\right| \\
& =\operatorname{dim} L\left(A^{1}\right)
\end{aligned}
$$

Hence $\operatorname{dim} L^{\sigma} \leq \operatorname{dim} L\left(A^{1}\right)$.
This shows that the homomorphism $L\left(A^{1}\right) \rightarrow L^{\sigma}$ is surjective and hence is an isomorphism. We note in particular that $\operatorname{dim}\left(L_{S}\right)^{\sigma}=1$ for each equivalence class $S$.

Thus we have shown that $L(A)^{\sigma}$ is isomorphic to $L\left(A^{1}\right)$. To be specific we have:

$$
\begin{aligned}
& L\left(A_{2 k}\right)^{\sigma} \cong L\left(B_{k}\right) \\
& L\left(A_{2 k-1}\right)^{\sigma} \cong L\left(C_{k}\right) \\
& L\left(D_{k+1}\right)^{\sigma} \cong L\left(B_{k}\right) \\
& L\left(D_{4}\right)^{\sigma} \cong L\left(G_{2}\right) \\
& L\left(E_{6}\right)^{\sigma} \cong L\left(F_{4}\right)
\end{aligned}
$$

where $\sigma$ is a graph automorphism of order $2,2,2,3,2$ respectively.

## 10

## Irreducible modules for semisimple Lie algebras

In the present chapter we shall determine the finite dimensional irreducible modules for a semisimple Lie algebra over $\mathbb{C}$. We begin by investigating certain important modules for such algebras known as Verma modules.

### 10.1 Verma modules

We begin with a lemma on universal enveloping algebras. Let $L$ be a finite dimensional Lie algebra over $\mathbb{C}$ and $K$ a subalgebra of $L$.

Lemma 10.1 There exists a unique algebra homomorphism $\theta: \mathfrak{U}(K) \rightarrow \mathfrak{H}(L)$ such that the diagram

$$
\begin{array}{cc}
K \xrightarrow{\sigma_{K}} \mathfrak{U}(K) \\
i \downarrow & \downarrow_{\theta} \\
L \xrightarrow[\sigma_{L}]{\rightarrow} \mathfrak{U}(L)
\end{array}
$$

commutes, where $i$ is the embedding of $K$ in $L$ and $\sigma_{K}, \sigma_{L}$ are the embeddings of $K, L$ in $\mathfrak{H}(K), \mathfrak{H}(L)$ respectively.

Also $\theta$ is injective.

Proof. Let $x \in K$. Then we must have

$$
\theta\left(\sigma_{K}(x)\right)=\sigma_{L}(i(x))
$$

Thus $\theta\left(\sigma_{K}(x)\right)$ is uniquely determined. Since $\mathfrak{H}(K)$ is generated by $\sigma_{K}(K)$ as algebra with 1 we see that $\theta$ is uniquely determined.

We now show that $\theta$ exists. We recall that

$$
\mathfrak{U}(L)=T(L) / J(L), \quad \mathfrak{U}(K)=T(K) / J(K)
$$

where $J(L)$ is the 2 -sided ideal of $T(L)$ generated by the elements

$$
x \otimes y-y \otimes x-[x y] \quad \text { for } x, y \in L
$$

Now the map $i: K \rightarrow L$ induces an algebra homomorphism $i: T(K) \rightarrow T(L)$ and we have

$$
i(x \otimes y-y \otimes x-[x y])=i(x) \otimes i(y)-i(y) \otimes i(x)-[i(x) i(y)]
$$

for all $x, y \in K$. This shows that

$$
i(J(K)) \subset J(L)
$$

Thus there is an algebra homomorphism $\theta: \mathfrak{U}(K) \rightarrow \mathfrak{U}(L)$ such that the required diagram commutes.

Finally we show that $\theta$ is injective. This follows from the PBW basis theorem 9.4. Let $x_{1}, \ldots, x_{r}$ be a basis of $K$. Suppose if possible there exists $u \in \mathfrak{U}(K)$ such that $u \neq 0$ and $\theta(u)=0$. Then $u$ is a non-zero linear combination of monomials $x_{1}^{e_{1}} \ldots x_{r}^{e_{r}}$. However, since $x_{1}, \ldots, x_{r}$ can be chosen as part of a basis of $L$, the PBW basis theorem for $L$ shows that such a combination of monomials cannot be zero in $L$. Hence $\theta(u) \neq 0$, a contradiction. Thus $\theta$ is injective.

This lemma shows that $\mathfrak{H}(K)$ may be regarded in a natural way as a subalgebra of $\mathfrak{H}(L)$.

We now suppose that $L$ is a finite dimensional semisimple Lie algebra over $\mathbb{C}$. Let $H$ be a Cartan subalgebra of $L$ and

$$
L=H \oplus \sum_{\alpha \in \Phi} L_{\alpha}
$$

be the Cartan decomposition of $L$ with respect to $H$. Let $\Phi^{+}$be the positive system of roots in $\Phi$. Then we have a triangular decomposition

$$
L=N^{-} \oplus H \oplus N
$$

where $N^{-}=\bigoplus_{\alpha \in \Phi^{-}} L_{\alpha}, N=\bigoplus_{\alpha \in \Phi^{+}} L_{\alpha}$. We recall that $H, N, N^{-}$are all subalgebras of $L$.

Let $B=H \oplus N$.

Lemma 10.2 (i) $B$ is a subalgebra of $L$.
(ii) $N$ is an ideal of $B$.
(iii) $B / N$ is isomorphic to $H$.

Proof. (i) We have

$$
[B B]=[H+N, H+N] \subset H+N=B
$$

since $H, N$ are subalgebras and $[H N] \subset N$.
(ii) $[N B]=[N, N+H] \subset N$.
(iii) $B / N=(H+N) / N \cong H / H \cap N \cong H$
since $H \cap N=0$, using Proposition 1.7.
Definition 10.3 Let $\lambda \in H^{*}$, i.e. $\lambda$ be a linear map from $H$ to $\mathbb{C}$. We recall that L has a basis

$$
\left\{e_{\alpha}, \alpha \in \Phi ; h_{i}, i=1, \ldots, l\right\}
$$

We define

$$
K_{\lambda}=\sum_{\alpha \in \Phi^{+}} \mathfrak{H}(L) e_{\alpha}+\sum_{i=1}^{l} \mathfrak{H}(L)\left(h_{i}-\lambda\left(h_{i}\right)\right)
$$

Thus $K_{\lambda}$ is the left ideal of $\mathfrak{H}(L)$ generated by the elements $e_{\alpha}, \alpha \in \Phi^{+}$, and $h_{i}-\lambda\left(h_{i}\right)$ for $i=1, \ldots, l$.
(We are as usual here embedding $L$ in $\mathfrak{1 l}(L)$.) We also define

$$
M(\lambda)=\mathfrak{U}(L) / K_{\lambda} .
$$

$M(\lambda)$ is a left $\mathfrak{H}(L)$-module called the Verma module determined by $\lambda$. It is our aim in this section to describe some of the properties of $M(\lambda)$.

We note that the elements $e_{\alpha}, \alpha \in \Phi^{+}$, and $h_{i}-\lambda\left(h_{i}\right)$ for $i=1, \ldots, l$ all lie in $\mathfrak{U}(B)$. We define

$$
K_{\lambda}^{\prime}=\sum_{\alpha \in \Phi^{+}} \mathfrak{H}(B) e_{\alpha}+\sum_{i=1}^{l} \mathfrak{H}(B)\left(h_{i}-\lambda\left(h_{i}\right)\right)
$$

to be the left ideal of $\mathfrak{H}(B)$ generated by these elements. Let

$$
\Phi^{+}=\left\{\beta_{1}, \ldots, \beta_{N}\right\}
$$

Then the set

$$
h_{1}, \ldots, h_{l}, \quad e_{\beta_{1}}, \ldots, e_{\beta_{N}}
$$

is a basis of $B$. It follows from the PBW basis theorem that the elements

$$
h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} \quad e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}} \quad s_{i} \geq 0 \quad t_{i} \geq 0
$$

form a basis of $\mathfrak{H}(B)$.

Proposition 10.4 (i) $\operatorname{dim} \mathfrak{U}(B) / K_{\lambda}^{\prime}=1$.
(ii) The elements

$$
\left(h_{1}-\lambda\left(h_{1}\right)\right)^{s_{1}} \ldots\left(h_{l}-\lambda\left(h_{l}\right)\right)^{s_{l}} e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}
$$

with $s_{i} \geq 0, t_{i} \geq 0$, excluding the element with $s_{i}=t_{i}=0$ for all $i$, form $a$ basis for $K_{\lambda}^{\prime}$.

Proof. It is not difficult to see that the elements

$$
\left(h_{1}-\lambda\left(h_{1}\right)\right)^{s_{1}} \ldots\left(h_{l}-\lambda\left(h_{l}\right)\right)^{s_{l}} e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}
$$

with $s_{i} \geq 0, t_{i} \geq 0$ also form a basis for $\mathfrak{H}(B)$. This can be seen, for example, by defining a partial ordering on the basis elements $h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}$. We say that $h_{1}^{s_{1}^{\prime}} \ldots h_{l}^{s_{l}^{s_{l}}} e_{\beta_{1}}^{t_{1}^{t_{1}}} \ldots e_{\beta_{N}}^{t_{N}^{\prime}}$ is lower than $h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}$ if $s_{1}^{\prime} \leq s_{1}, \ldots$, $s_{l}^{\prime} \leq s_{l}, \quad t_{1}^{\prime}=t_{1}, \ldots, t_{N}^{\prime}=t_{N}$. Then there are only a finite number of basis elements lower than a given one, and the element

$$
\left(h_{1}-\lambda\left(h_{1}\right)\right)^{s_{1}} \ldots\left(h_{l}-\lambda\left(h_{l}\right)\right)^{s_{l}} e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}
$$

is the sum of $h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}$ with a linear combination of strictly lower basis elements in the partial order. An induction argument on the partial order will then show that the elements

$$
\left(h_{1}-\lambda\left(h_{1}\right)\right)^{s_{1}} \ldots\left(h_{l}-\lambda\left(h_{l}\right)\right)^{s_{l}} e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}
$$

$s_{i} \geq 0, t_{i} \geq 0$ span $\mathfrak{U}(B)$ and are linearly independent. Now all these elements clearly lie in $K_{\lambda}^{\prime}$, with the exception of the element with $s_{i}=0, t_{i}=0$ for all $i$. This is the unit element 1 . However, 1 does not lie in $K_{\lambda}^{\prime}$, as the following argument shows.

Consider the representation $\lambda$ of $H$ mapping $h_{i}$ to $\lambda\left(h_{i}\right)$ for $i=1, \ldots, l$. Since $B / N$ is isomorphic to $H$ there is a 1-dimensional representation $\lambda$ of $B$ with $N$ in the kernel agreeing with the above representation on $B / N \cong H$. This in turn gives a 1-dimensional representation $\rho$ of $\mathfrak{H}(B)$ under which

$$
\begin{aligned}
e_{\alpha} & \rightarrow 0 \quad \alpha \in \Phi^{+} \\
h_{i} & \rightarrow \lambda\left(h_{i}\right) \quad i=1, \ldots, l \\
1 & \rightarrow 1
\end{aligned}
$$

Now ker $\rho$ is a 2-sided ideal of $\mathfrak{H}(B)$ containing $e_{\alpha}, \alpha \in \Phi^{+}$and $h_{i}-\lambda\left(h_{i}\right)$ so containing $K_{\lambda}^{\prime}$. Thus we have

$$
K_{\lambda}^{\prime} \subset \operatorname{ker} \rho \quad 1 \notin \operatorname{ker} \rho
$$

hence $1 \notin K_{\lambda}^{\prime}$.

It follows that the elements

$$
\left(h_{1}-\lambda\left(h_{1}\right)\right)^{s_{1}} \ldots\left(h_{l}-\lambda\left(h_{l}\right)\right)^{s_{l}} e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}
$$

for $s_{i} \geq 0, t_{i} \geq 0$, excluding 1 , form a basis of $K_{\lambda}^{\prime}$ and that

$$
\operatorname{dim} \mathfrak{U}(B) / K_{\lambda}^{\prime}=1
$$

Proposition $10.5 K_{\lambda} \cap \mathfrak{U}\left(N^{-}\right)=O$.

Proof. We are here regarding $\mathfrak{H}\left(N^{-}\right)$as a subalgebra of $\mathfrak{H}(L)$ as in Lemma 10.1. Now we have

$$
L=N^{-} \oplus B .
$$

Regarding $\mathfrak{H}(B)$ as a subalgebra of $\mathfrak{H}(L)$ also we assert that

$$
\mathfrak{U}(L)=\mathfrak{U l}\left(N^{-}\right) \mathfrak{U}(B),
$$

i.e. each element of $\mathfrak{U}(L)$ is a finite sum $\sum x_{i} y_{i}$ with $x_{i} \in \mathfrak{U}\left(N^{-}\right)$, $y_{i} \in \mathfrak{U}(B)$. This follows from the PBW basis theorem, choosing bases for $N^{-}$and for $B$ and combining them to give a basis of $L$. We then have

$$
\begin{aligned}
K_{\lambda} & =\sum_{\alpha \in \Phi^{+}} \mathfrak{H}(L) e_{\alpha}+\sum_{i=1}^{l} \mathfrak{H}(L)\left(h_{i}-\lambda\left(h_{i}\right)\right) \\
& =\sum_{\alpha \in \Phi^{+}} \mathfrak{H}\left(N^{-}\right) \mathfrak{U}(B) e_{\alpha}+\sum_{i=1}^{l} \mathfrak{H}\left(N^{-}\right) \mathfrak{U}(B)\left(h_{i}-\lambda\left(h_{i}\right)\right) \\
& =\mathfrak{U}\left(N^{-}\right)\left(\sum_{\alpha \in \Phi^{+}} \mathfrak{U}(B) e_{\alpha}+\sum_{i=1}^{l} \mathfrak{U}(B)\left(h_{i}-\lambda\left(h_{i}\right)\right)\right)=\mathfrak{U}\left(N^{-}\right) K_{\lambda}^{\prime} .
\end{aligned}
$$

It follows that each element of $K_{\lambda}$ is a linear combination of terms of form

$$
f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}}\left(h_{1}-\lambda\left(h_{1}\right)\right)^{s_{1}} \ldots\left(h_{l}-\lambda\left(h_{l}\right)\right)^{s_{l}} e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}
$$

where $f_{\alpha}=e_{-\alpha}, \alpha \in \Phi^{+}$, and $r_{i} \geq 0, s_{i} \geq 0, t_{i} \geq 0$ with $\left(s_{1}, \ldots s_{l}, t_{1}, \ldots t_{N}\right) \neq$ $(0, \ldots 0,0 \ldots 0)$. No non-zero element of $\mathfrak{U}\left(N^{-}\right)$can be a linear combination of such terms by the PBW basis theorem. Hence

$$
K_{\lambda} \cap \mathfrak{U}\left(N^{-}\right)=O
$$

Let $m_{\lambda} \in M(\lambda)$ be defined by $m_{\lambda}=1+K_{\lambda}$. Thus 1 maps to $m_{\lambda}$ under the natural homomorphism

$$
\mathfrak{U}(L) \rightarrow \mathfrak{U}(L) / K_{\lambda}=M(\lambda)
$$

Theorem 10.6 (i) Each element of $M(\lambda)$ is uniquely expressible in the form um ${ }_{\lambda}$ for some $u \in \mathfrak{U}\left(N^{-}\right)$.
(ii) The elements $f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda}$ for all $r_{i} \geq 0$ form a basis for $M(\lambda)$.

Proof. Each element of $\mathfrak{H}(L)$ has form $u \cdot 1$ for some $u \in \mathfrak{H}(L)$. Thus each element of $M(\lambda)=\mathfrak{U}(L) / K_{\lambda}$ has form $u m_{\lambda}$ for some $u \in \mathfrak{U}(L)$.

Now $f_{\beta_{1}}, \ldots, f_{\beta_{N}}, h_{1}, \ldots, h_{l}, \quad e_{\beta_{1}}, \ldots, e_{\beta_{N}}$ are a basis of $L$ so the elements

$$
f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} \quad h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} \quad e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}
$$

$r_{i} \geq 0, s_{i} \geq 0, t_{i} \geq 0$ form a basis of $\mathfrak{H}(L)$ by the PBW basis theorem. Thus $u$ is a linear combination of such elements, and $u m_{\lambda}$ is a linear combination of elements

$$
f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} \quad h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} \quad e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}} m_{\lambda}
$$

Now this element is 0 if any $t_{i}$ is positive. Suppose then that all $t_{i}=0$. Then

$$
h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} m_{\lambda}=\gamma m_{\lambda} \quad \text { for some } \gamma \in \mathbb{C}
$$

since $h_{i} m_{\lambda}=\lambda\left(h_{i}\right) m_{\lambda}$. Thus $u m_{\lambda}$ is a linear combination of elements of form

$$
f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda}
$$

Thus elements of this form for $r_{i} \geq 0$ span $M(\lambda)$. They are also linearly independent. For if we have

$$
\sum_{r_{1}, \ldots, r_{N}} \xi_{r_{1}, \ldots, r_{N}} f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda}=0
$$

with $\xi_{r_{1}, \ldots, r_{N}} \in \mathbb{C}$ then it follows that

$$
\sum_{r_{1}, \ldots, r_{N}} \xi_{r_{1}, \ldots, r_{N}} f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} \in K_{\lambda} \cap \mathfrak{U}\left(N^{-}\right)
$$

Hence this element is 0 by Proposition 10.5. Thus each $\xi_{r_{1}, \ldots, r_{N}}=0$ by the PBW basis theorem for $\mathfrak{U}\left(N^{-}\right)$.

Thus the elements $f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda}$ for $r_{1} \geq 0, \ldots, r_{N} \geq 0$ form a basis for $M(\lambda)$. It follows that each element of $M(\lambda)$ is uniquely expressible in the form $u m_{\lambda}$ for $u \in \mathfrak{U}\left(N^{-}\right)$.

We now regard $M(\lambda)$ as an $H$-module. For each 1-dimensional representation $\mu$ of $H$ we define

$$
M(\lambda)_{\mu}=\{m \in M(\lambda) ; x m=\mu(x) m \quad \text { for all } x \in H\} .
$$

$M(\lambda)_{\mu}$ is a subspace of $M(\lambda)$.

Theorem 10.7 (i) $M(\lambda)=\bigoplus_{\mu \in H^{*}} M(\lambda)_{\mu}$.
(ii) $M(\lambda)_{\mu} \neq 0$ if and only if $\lambda-\mu$ is a sum of positive roots.
(iii) $\operatorname{dim} M(\lambda)_{\mu}=\mathfrak{B}(\lambda-\mu)$, the number of ways of expressing $\lambda-\mu$ as a sum of positive roots.
$\mathfrak{B}(\lambda-\mu)$ is the number of vectors $\left(r_{1}, \ldots, r_{N}\right)$ with $r_{i} \in \mathbb{Z}, r_{i} \geq 0$ such that

$$
\lambda-\mu=r_{1} \beta_{1}+\cdots+r_{N} \beta_{N} .
$$

Proof. We know from Theorem 10.6 that the elements $f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda}$ with $r_{i} \geq 0$ form a basis for $M(\lambda)$. We show that

$$
x f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda}=\left(\lambda-r_{1} \beta_{1}-\cdots-r_{N} \beta_{N}\right)(x) f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda} \quad \text { for all } x \in H
$$

We prove this by induction on $r_{1}+\cdots+r_{N}$, the result being clear if all $r_{i}=0$. So suppose not all $r_{i}$ are 0 and let $i$ be the least integer with $r_{i}>0$. Then we have

$$
x f_{\beta_{i}}^{r_{i}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda}=f_{\beta_{i}} x f_{\beta_{i}}^{r_{i}-1} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda}-\beta_{i}(x) f_{\beta_{i}}^{r_{i}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda}
$$

It follows that

$$
x f_{\beta_{i}}^{r_{i}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda}=\left(\lambda-r_{i} \beta_{i}-\cdots-r_{N} \beta_{N}\right)(x) f_{\beta_{i}}^{r_{i}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda}
$$

as required.
This implies that $f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda} \in M(\lambda)_{\mu}$ where $\mu=\lambda-r_{1} \beta_{1}-\cdots-r_{N} \beta_{N}$. Since these elements form a basis of $M(\lambda)$ we see that

$$
M(\lambda)=\sum_{\mu} M(\lambda)_{\mu}
$$

We now show this sum is direct. To see this we must show that if a finite sum $\sum v_{\mu}$ is 0 with $v_{\mu} \in M(\lambda)_{\mu}$, then each $v_{\mu}$ is 0 .

It is sufficient to show that

$$
M(\lambda)_{\mu} \cap\left(M(\lambda)_{\mu_{1}}+\cdots+M(\lambda)_{\mu_{k}}\right)=O
$$

where the elements $\mu, \mu_{1}, \ldots, \mu_{k} \in H^{*}$ are all distinct.
Let $v$ lie in this intersection. Then we have

$$
v=v_{\mu_{1}}+\cdots+v_{\mu_{k}}
$$

where $v \in M(\lambda)_{\mu}, \quad v_{\mu_{i}} \in M(\lambda)_{\mu_{i}}$. Thus

$$
\begin{aligned}
& (x-\mu(x)) v=0 \\
& \left(x-\mu_{i}(x)\right) v_{\mu_{i}}=0
\end{aligned}
$$

for all $x \in H$. Hence

$$
\left(x-\mu_{1}(x)\right) \ldots\left(x-\mu_{k}(x)\right)\left(v_{\mu_{1}}+\cdots+v_{\mu_{k}}\right)=0
$$

that is $\left(x-\mu_{1}(x)\right) \ldots\left(x-\mu_{k}(x)\right) v=0$. Since the vector space $H$ over $\mathbb{C}$ cannot be expressed as the union of finitely many proper subspaces we can find $x \in H$ such that

$$
\mu(x) \neq \mu_{1}(x), \quad \mu(x) \neq \mu_{2}(x), \quad \ldots, \quad \mu(x) \neq \mu_{k}(x)
$$

Thus the polynomials

$$
t-\mu(x), \quad\left(t-\mu_{1}(x)\right)\left(t-\mu_{2}(x)\right) \ldots\left(t-\mu_{k}(x)\right)
$$

in $\mathbb{C}[t]$ for this element $x$ are coprime. Thus there exist polynomials $p(t), q(t) \in \mathbb{C}[t]$ such that

$$
p(t)(t-\mu(x))+q(t)\left(t-\mu_{1}(x)\right) \ldots\left(t-\mu_{k}(x)\right)=1 .
$$

Thus we have

$$
p(x)(x-\mu(x))+q(x)\left(x-\mu_{1}(x)\right) \ldots\left(x-\mu_{k}(x)\right)=1 .
$$

It follows that

$$
p(x)(x-\mu(x)) v+q(x)\left(x-\mu_{1}(x)\right) \ldots\left(x-\mu_{k}(x)\right) v=v .
$$

The above conditions show that the left-hand side is zero, hence we have $v=0$. Thus

$$
M(\lambda)=\bigoplus_{\mu} M(\lambda)_{\mu}
$$

Let $\Lambda=\left\{\mu \in H^{*} ; \lambda-\mu\right.$ is a sum of positive roots $\}$. For each $\lambda \in \Lambda$ let $N_{\mu}$ be the subspace of $M(\lambda)$ spanned by the basis vectors

$$
f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} m_{\lambda}
$$

with $\lambda-r_{1} \beta_{1}-\cdots-r_{N} \beta_{N}=\mu$. Since these vectors for all such $\mu$ form a basis of $M(\lambda)$ we have

$$
M(\lambda)=\bigoplus_{\mu \in \Lambda} N_{\mu}
$$

On the other hand we know that

$$
N_{\mu} \subset M(\lambda)_{\mu}
$$

and $M(\lambda)=\bigoplus_{\mu \in H^{*}} M(\lambda)_{\mu}$. It follows that $M(\lambda)_{\mu}=N_{\mu}$ for all $\mu \in \Lambda$, and that $M(\lambda)_{\mu}=O$ for all $\mu \in H^{*}$ with $\mu \notin \Lambda$.

Thus we have $\operatorname{dim} M(\lambda)_{\mu}=\operatorname{dim} N_{\mu}=$ the number of vectors $\left(r_{1}, \ldots, r_{N}\right)$ with $r_{i} \in \mathbb{Z}, r_{i} \geq 0$ such that

$$
\lambda-r_{1} \beta_{1}-\cdots-r_{N} \beta_{N}=\mu
$$

This gives $\operatorname{dim} M(\lambda)_{\mu}=\mathfrak{B}(\lambda-\mu)$ as required.

Definition $10.8 \mu \in H^{*}$ is called a weight of $M(\lambda)$ if $M(\lambda)_{\mu} \neq O$, and $M(\lambda)_{\mu}$ is called the weight space of $M(\lambda)$ with weight $\mu$.

We note that since $M(\lambda)=\oplus_{\mu} M(\lambda)_{\mu}$ an element $m$ of $M(\lambda)$ satisfying the condition that, for all $x \in H,(x-\mu(x))^{k} m=0$ for some $k>0$ can have no non-zero component in any $M(\lambda)_{\nu}$ for $\nu \neq \mu$, and must therefore lie in $M(\lambda)_{\mu}$. Thus we have

$$
\begin{aligned}
M(\lambda)_{\mu}=\{ & m \in M(\lambda) ; \text { for each } x \in H \text { there exists } k \text { such that } \\
& \left.(x-\mu(x))^{k} m=0\right\} .
\end{aligned}
$$

This shows that our definitions of weight and weight space here in the context of H -modules are compatible with the definitions in Chapter 2 in the context of representations of nilpotent Lie algebras.

Theorem 10.7 asserts that a Verma module is the direct sum of its weight spaces. There are infinitely many weights, but each weight space is finite dimensional.

We now proceed to another very important property of Verma modules.

Theorem 10.9 $M(\lambda)$ has a unique maximal submodule.

Proof. Let $V$ be a $\mathfrak{H}(L)$-submodule of $M(\lambda)$ with $V \neq M(\lambda)$. Let $v \in V$. By Theorem 10.7 we have

$$
v=\sum_{i} v_{\mu_{i}} \quad v_{\mu_{i}} \in M(\lambda)_{\mu_{i}}
$$

summed over a finite set of distinct weights $\mu_{i}$. We aim to show that each $v_{\mu_{i}}$ lies in $V$ also. We have

$$
x v_{\mu_{i}}=\mu_{i}(x) v_{\mu_{i}} \quad x \in H .
$$

Hence

$$
\prod_{\substack{j \\ j \neq i}}\left(x-\mu_{j}(x)\right) v=\prod_{\substack{j \\ j \neq i}}\left(x-\mu_{j}(x)\right) v_{\mu_{i}}=\prod_{\substack{j \\ j \neq i}}\left(\mu_{i}(x)-\mu_{j}(x)\right) v_{\mu_{i}} .
$$

Since $H$ is not the union of finitely many proper subspaces we can find $x \in H$ with $\mu_{i}(x) \neq \mu_{j}(x)$ for all $j \neq i$. For such an $x$ we have

$$
\prod_{\substack{j \\ j \neq i}}\left(\mu_{i}(x)-\mu_{j}(x)\right) v_{\mu_{i}} \in V
$$

and

$$
\prod_{\substack{j \\ j \neq i}}\left(\mu_{i}(x)-\mu_{j}(x)\right) \neq 0
$$

It follows that $v_{\mu_{i}} \in V$.
We now define $V_{\mu}=V \cap M(\lambda)_{\mu}$. We have shown that $V=\sum_{\mu} V_{\mu}$. Since we know that $M(\lambda)=\bigoplus_{\mu} M(\lambda)_{\mu}$ it follows that the sum in $V$ must be direct, that is

$$
V=\bigoplus_{\mu} V_{\mu}
$$

Thus every submodule $V$ of $M(\lambda)$ is also the direct sum of its weight spaces. Now $V_{\lambda}=O$. For if $V_{\lambda} \neq O$ then $V_{\lambda}=M(\lambda)_{\lambda}$ since $\operatorname{dim} M(\lambda)_{\lambda}=1$. This would imply that $m_{\lambda} \in V$. But then

$$
M(\lambda)=\mathfrak{H}(L) m_{\lambda} \subset V
$$

so $V=M(\lambda)$, a contradiction. Thus $V_{\lambda}=O$ and we have

$$
V=\bigoplus_{\substack{\mu \\ \mu \neq \lambda}} V_{\mu} \subset \sum_{\substack{\mu \\ \mu \neq \lambda}} M(\lambda)_{\mu}
$$

Thus every proper submodule $V$ of $M(\lambda)$ lies in the subspace $\sum_{\mu \neq \lambda}^{\mu} M(\lambda)_{\mu}$ of codimension 1 in $M(\lambda)$. Let $J(\lambda)$ be the sum of all the proper submodules of $M(\lambda) . J(\lambda)$ lies in the above subspace of codimension 1 , so is a proper submodule of $M(\lambda)$. Thus $J(\lambda)$ is the unique maximal submodule of $M(\lambda)$, since it contains all proper submodules of $M(\lambda)$.

Definition 10.10 Let $\lambda, \mu \in H^{*}$. In view of Theorem 10.7 it is natural to make the following definition.

We say that $\lambda \succ \mu$ if $\lambda-\mu$ is a sum of positive roots. This is a partial order on $H^{*}$.

Theorem 10.7 shows that the weights of $M(\lambda)$ are precisely the $\mu \in H^{*}$ with $\mu \prec \lambda$. Thus $\lambda$ is the highest weight of $M(\lambda)$ with respect to this partial order.
$M(\lambda)$ is called the Verma module with highest weight $\lambda$.

We also define $L(\lambda)=M(\lambda) / J(\lambda)$. Since $J(\lambda)$ is a maximal submodule of $M(\lambda), L(\lambda)$ is an irreducible $\mathfrak{H}(L)$-module. In subsequent sections of this chapter we shall determine under what circumstances the irreducible module $L(\lambda)$ is finite dimensional. We note that $\lambda$ is a weight of $L(\lambda)$, since $J(\lambda)_{\lambda}=$ $O$. Thus $\operatorname{dim} L(\lambda)_{\lambda}=1$ and $\lambda$ is the highest weight of $L(\lambda)$.

### 10.2 Finite dimensional irreducible modules

Now let $V$ be any finite dimensional irreducible $L$-module where, as usual in this chapter, $L$ is a finite dimensional semisimple Lie algebra over $\mathbb{C}$. Let $H$ be a Cartan subalgebra of $L$ and

$$
\left\{e_{\alpha}, \alpha \in \Phi^{+} ; h_{i}, i=1, \ldots, l ; f_{\alpha}, \alpha \in \Phi^{+}\right\}
$$

be a basis of $L$ adapted to $H$. We may regard $V$ as an $H$-module. Now $H$ is abelian, so in particular nilpotent, thus we may apply the representation theory of nilpotent Lie algebras developed in Chapter 2. By Theorem 2.9 we have

$$
V=\bigoplus_{\lambda} V_{\lambda}
$$

where $V_{\lambda}=\left\{v \in V\right.$; for each $x \in H$ there exists $k$ such that $\left.(x-\lambda(x))^{k} v=0\right\}$. We also know from Chapter 2 that each non-zero $V_{\lambda}$ contains a non-zero vector $v$ such that

$$
x v=\lambda(x) v \quad \text { for all } x \in H
$$

We shall show that in our present situation the weight spaces $V_{\lambda}$ can be defined more simply.

Proposition 10.11 Let $W_{\lambda}=\{v \in V ; x v=\lambda(x) v$ for all $x \in H\}$. Then $W_{\lambda}=V_{\lambda}$.

Proof. It is clear that $W_{\lambda} \subset V_{\lambda}$ and that $W_{\lambda} \neq O$ whenever $V_{\lambda} \neq O$. Let $W=$ $\sum_{\lambda} W_{\lambda}$. Since $V=\bigoplus_{\lambda} V_{\lambda}$ and $W_{\lambda} \subset V_{\lambda}$ we see that $W=\bigoplus_{\lambda} W_{\lambda}$. We shall show that $W$ is a submodule of $V$. To see this it is sufficient to show that $h_{i} w, e_{\alpha} w, f_{\alpha} w$ lie in $W$ for all $w \in W_{\lambda}$, all $i=1, \ldots, l$ and all $\alpha \in \Phi^{+}$. Now we have

$$
\begin{aligned}
& h_{i} w=\lambda\left(h_{i}\right) w \in W \\
& \begin{aligned}
x\left(e_{\alpha} w\right) & =e_{\alpha}(x w)+\alpha(x) e_{\alpha} w \\
& =\lambda(x) e_{\alpha} w+\alpha(x) e_{\alpha} w \\
& =(\lambda+\alpha)(x) e_{\alpha} w .
\end{aligned}
\end{aligned}
$$

Hence $e_{\alpha} w \in W_{\lambda+\alpha} \subset W$. Similarly we have $f_{\alpha} w \in W_{\lambda-\alpha} \subset W$. Thus $W$ is a $\mathfrak{U}(L)$-submodule of $V$. Since $W \neq O$ and $V$ is irreducible we have $W=V$. It follows that $W_{\lambda}=V_{\lambda}$ for each $\lambda \in H^{*}$.

Thus the irreducible module $V$ is the direct sum of its weight spaces $V_{\lambda}$ and $V_{\lambda}$ is the set of all $v \in V$ such that $x v=\lambda(x) v$ for all $x \in H$.

We now consider the set of all weights $\lambda$ for $V$, that is the set of all $\lambda \in H^{*}$ for which $V_{\lambda} \neq O$. This is a finite set, so will contain at least one weight maximal in the partial order $\succ$ defined in Definition 10.10. Let $\lambda$ be such a weight of $V$. If $\mu \succ \lambda$ and $\mu \neq \lambda$ then $\mu$ is not a weight of $V$.

We may choose $v_{\lambda} \in V_{\lambda}$ with $v_{\lambda} \neq 0$.

Proposition 10.12 (i) $x v_{\lambda}=\lambda(x) v_{\lambda}$ for all $x \in H$.
(ii) $e_{\alpha} v_{\lambda}=0$ for all $\alpha \in \Phi^{+}$.
(iii) $V=\mathfrak{H}\left(N^{-}\right) v_{\lambda}$
(iv) $\lambda$ is the highest weight of $V$.

Proof. Condition (i) is clear. We have

$$
x\left(e_{\alpha} v_{\lambda}\right)=e_{\alpha}\left(x v_{\lambda}\right)+\alpha(x) e_{\alpha} v_{\lambda}
$$

for all $x \in H$. Now if $e_{\alpha} v_{\lambda} \neq 0$ this implies that $\lambda+\alpha$ is a weight of $V$. But $\lambda+\alpha \succ \lambda$ so this cannot be the case. Hence $e_{\alpha} v_{\lambda}=0$ for $\alpha \in \Phi^{+}$.

Now $V=\mathfrak{U}(L) v_{\lambda}$ since $v_{\lambda} \neq 0$ and $V$ is an irreducible $\mathfrak{U}(L)$-module. Thus each element of $V$ is a linear combination of elements of the form

$$
f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} \quad h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} \quad e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}} v_{\lambda} .
$$

This element is 0 unless all $t_{i}$ are 0 . In that case it is a scalar multiple of

$$
f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} v_{\lambda}
$$

Hence $V=\mathfrak{U}\left(N^{-}\right) v_{\lambda}$.
Finally we have

$$
x\left(f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} v_{\lambda}\right)=\left(\lambda-r_{1} \beta_{1}-\cdots-r_{N} \beta_{N}\right)(x) f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} v_{\lambda}
$$

as in the proof of Theorem 10.7 ; thus all weights of $V$ have form

$$
\mu=\lambda-r_{1} \beta_{1}-\cdots-r_{N} \beta_{N}
$$

Thus $\lambda \succ \mu$ for all weights $\mu$ of $V$.

It follows from Proposition 10.12 that the set of weights of $V$ has a unique maximal element $\lambda$ with respect to the partial order $\succ$.

We now compare the finite dimensional module $V$ with the Verma module $M(\lambda)$.

Proposition 10.13 There exists a surjective homomorphism $\theta: M(\lambda) \rightarrow V$ of $\mathfrak{U}(L)$-modules such that $\theta\left(m_{\lambda}\right)=v_{\lambda}$.

Proof. We recall from Theorem 10.6 that each element of $M(\lambda)$ is uniquely expressible in the form $u m_{\lambda}$ with $u \in \mathfrak{U}\left(N^{-}\right)$. We define a linear map

$$
\theta: M(\lambda) \rightarrow V
$$

by $\theta\left(u m_{\lambda}\right)=u v_{\lambda} \quad u \in \mathfrak{H}\left(N^{-}\right)$. Then $\theta$ is surjective by Proposition 10.12 (iii). We must check that $\theta$ is a homomorphism of $\mathfrak{U}(L)$-modules. Thus we must show

$$
\theta\left(\text { yum }_{\lambda}\right)=y u v_{\lambda} \quad \text { for all } y \in \mathfrak{H}(L)
$$

By the PBW basis theorem we know that the element $y u$ of $\mathfrak{U}(L)$ can be written as a finite sum

$$
y u=\sum_{i} a_{i} b_{i} c_{i}
$$

where $a_{i} \in \mathfrak{U}\left(N^{-}\right), b_{i} \in \mathfrak{H}(H), c_{i} \in \mathfrak{H}(N)$. Thus

$$
\text { yum }_{\lambda}=\sum_{i} a_{i} b_{i} c_{i} m_{\lambda}
$$

Now $b_{i} c_{i} m_{\lambda}=\xi_{i} m_{\lambda}$ for some $\xi_{i} \in \mathbb{C}$. Hence

$$
y u m_{\lambda}=\left(\sum_{i} \xi_{i} a_{i}\right) m_{\lambda}
$$

Since $\sum_{i} \xi_{i} a_{i} \in \mathfrak{H}\left(N^{-}\right)$we have

$$
\theta\left(y u m_{\lambda}\right)=\theta\left(\left(\sum_{i} \xi_{i} a_{i}\right) m_{\lambda}\right)=\left(\sum_{i} \xi_{i} a_{i}\right) v_{\lambda}
$$

On the other hand we have

$$
y u v_{\lambda}=\left(\sum_{i} a_{i} b_{i} c_{i}\right) v_{\lambda}=\left(\sum_{i} \xi_{i} a_{i}\right) v_{\lambda}
$$

since $b_{i} c_{i} v_{\lambda}=\xi_{i} v_{\lambda}$. Hence $\theta\left(y u m_{\lambda}\right)=y u v_{\lambda}$ for all $y \in \mathfrak{H}(L)$. Thus $\theta$ is a homomorphism of $\mathfrak{U}(L)$-modules.

Corollary 10.14 $V$ is isomorphic to $L(\lambda)$.

Proof. Since $V$ is irreducible the kernel of $\theta$ must be a maximal submodule of $M(\lambda)$. But $M(\lambda)$ has a unique maximal submodule $J(\lambda)$, by Theorem 10.9. Thus ker $\theta=J(\lambda)$. Hence $V$ is isomorphic to $M(\lambda) / J(\lambda)=L(\lambda)$.

Thus we have seen that every finite dimensional irreducible $L$-module is isomorphic to one of the irreducible modules $L(\lambda)$ obtained as irreducible quotients of Verma modules. However, we shall see that by no means all the $L(\lambda)$ are finite dimensional.

Proposition 10.15 Suppose $L(\lambda)$ is finite dimensional. Then $\lambda\left(h_{i}\right)$ is a nonnegative integer for each $i=1, \ldots, l$.

Proof. Let $v_{\lambda}$ be a highest weight vector of $L(\lambda)$, that is $v_{\lambda} \in L(\lambda)_{\lambda}$ and $v_{\lambda} \neq 0$. As in Section 7.1 we shall choose elements $e_{i} \in L_{\alpha_{i}}, f_{i} \in L_{-\alpha_{i}}$ such that $\left[e_{i} f_{i}\right]=h_{i}$. We consider the sequence of elements

$$
v_{\lambda}, \quad f_{i} v_{\lambda}, \quad f_{i}^{2} v_{\lambda}, \quad \ldots
$$

of $L(\lambda)$. We have

$$
x\left(f_{i}^{k} v_{\lambda}\right)=\left(\lambda-k \alpha_{i}\right)(x)\left(f_{i}^{k} v_{\lambda}\right)
$$

for all $x \in H$. Thus we have

$$
v_{\lambda} \in L(\lambda)_{\lambda}, \quad f_{i} v_{\lambda} \in L(\lambda)_{\lambda-\alpha_{i}}, \quad f_{i}^{2} v_{\lambda} \in L(\lambda)_{\lambda-2 \alpha_{i}}
$$

and so on. Now $L(\lambda)$, being finite dimensional, has only finitely many distinct weights. Thus there exists $p \in \mathbb{Z}, p \geq 0$ such that

$$
\begin{aligned}
& f_{i}^{k} v_{\lambda} \neq 0 \quad \text { for } k \leq p \\
& f_{i}^{p+1} v_{\lambda}=0
\end{aligned}
$$

Let $M=\mathbb{C} v_{\lambda}+\mathbb{C} f_{i} v_{\lambda}+\cdots+\mathbb{C} f_{i}^{p} v_{\lambda}$. This sum is direct since $v_{\lambda}, f_{i} v_{\lambda}, \ldots, f_{i}^{p} v_{\lambda}$ all lie in different weight spaces. We show that $M$ is a submodule with respect to the subalgebra $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$ of $L$. It is clear from the definitions that $h_{i} M \subset M$ and $f_{i} M \subset M$. We shall show that $e_{i} M \subset M$ also.

We verify that $e_{i} f_{i}^{k} v_{\lambda} \in M$ by induction on $k$. If $k=0$ we have $e_{i} v_{\lambda}=0$. If $k>0$ we have

$$
e_{i} f_{i}^{k} v_{\lambda}=f_{i} e_{i} f_{i}^{k-1} v_{\lambda}+h_{i} f_{i}^{k-1} v_{\lambda}
$$

Now $e_{i} f_{i}^{k-1} v_{\lambda} \in M$ by induction, hence $e_{i} f_{i}^{k} v_{\lambda} \in M$ also. Thus $M$ is an $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$-submodule. We consider the trace of $h_{i}$ on $M$. This can be calculated in two ways. On the one hand we have

$$
\operatorname{trace}_{M} h_{i}=\operatorname{trace}_{M}\left[e_{i} f_{i}\right]=\operatorname{trace}_{M}\left(e_{i} f_{i}-f_{i} e_{i}\right)=0 .
$$

On the other hand we have

$$
\begin{aligned}
\operatorname{trace}_{M} h_{i} & =\lambda\left(h_{i}\right)+\left(\lambda-\alpha_{i}\right)\left(h_{i}\right)+\cdots+\left(\lambda-p \alpha_{i}\right)\left(h_{i}\right) \\
& =(p+1) \lambda\left(h_{i}\right)-p(p+1)
\end{aligned}
$$

since $\alpha_{i}\left(h_{i}\right)=2$. Hence

$$
\operatorname{trace}_{M} h_{i}=(p+1) \lambda\left(h_{i}\right)-p(p+1) .
$$

It follows that

$$
(p+1)\left(\lambda\left(h_{i}\right)-p\right)=0,
$$

that is $\lambda\left(h_{i}\right)=p$. Thus $\lambda\left(h_{i}\right) \in \mathbb{Z}$ and $\lambda\left(h_{i}\right) \geq 0$.
The condition $\lambda\left(h_{i}\right) \in \mathbb{Z}, \lambda\left(h_{i}\right) \geq 0$ for all $i=1, \ldots, l$ is therefore necessary for $L(\lambda)$ to be finite dimensional. In the next section we shall show that this condition is also sufficient.

### 10.3 The finite dimensionality criterion

We consider the set of $\lambda \in H^{*}$ such that $\lambda\left(h_{i}\right) \in \mathbb{Z}, \lambda\left(h_{i}\right) \geq 0$ for $i=1, \ldots, l$.
Definition 10.16 Let $\omega_{i} \in H^{*}$ be the element satisfying $\omega_{i}\left(h_{i}\right)=1, \omega_{i}\left(h_{j}\right)=0$ if $j \neq i$. The elements $\omega_{1}, \ldots, \omega_{l} \in H^{*}$ are called the fundamental weights.

We note that $\omega_{1}, \ldots, \omega_{l}$ are linearly independent, since this is true of $h_{1}, \ldots, h_{l} \in H$. Thus $\omega_{1}, \ldots, \omega_{l}$ form a basis of $H^{*}$. Let $X=\left\{n_{1} \omega_{1}+\cdots\right.$ $\left.+n_{l} \omega_{l} ; n_{1}, \ldots, n_{l} \in \mathbb{Z}\right\} . X$ is a free abelian subgroup of $H^{*}$ with basis the set of fundamental weights and is called the lattice of integral weights or, briefly, the weight lattice. It is clear that an element $\lambda \in H^{*}$ lies in $X$ if and only if $\lambda\left(h_{i}\right) \in \mathbb{Z}$ for $i=1, \ldots, l$. Let $X^{+}=\left\{n_{1} \omega_{1}+\cdots\right.$ $+n_{l} \omega_{l} ; n_{i} \in \mathbb{Z}, n_{i} \geq 0$ for $\left.i=1, \ldots, l\right\} . X^{+}$is called the set of dominant integral weights.

An element $\lambda \in H^{*}$ lies in $X^{+}$if and only if $\lambda\left(h_{i}\right) \in \mathbb{Z}$ and $\lambda\left(h_{i}\right) \geq 0$ for $i=1, \ldots, l$.

We have seen, therefore, that if $L(\lambda)$ is finite dimensional then $\lambda$ is a dominant integral weight, and wish to prove the converse.

We shall first explore the connection between the fundamental weights $\omega_{1}, \ldots, \omega_{l}$ and the fundamental roots $\alpha_{1}, \ldots, \alpha_{l}$.

Proposition $10.17 \alpha_{i}=\sum_{j} A_{j i} \omega_{j}$. Thus the matrix expressing the fundamental roots as linear combinations of the fundamental weights is the transpose of the Cartan matrix.

Proof. Since $\omega_{1}, \ldots, \omega_{l}$ are a basis for $H^{*}$ there exist $c_{i j} \in \mathbb{C}$ such that

$$
\alpha_{i}=\sum_{j} c_{i j} \omega_{j}
$$

Then we have

$$
\alpha_{i}\left(h_{j}\right)=c_{i j}
$$

Hence

$$
c_{i j}=\alpha_{i}\left(h_{j}\right)=\alpha_{i}\left(\frac{2 h_{\alpha_{j}}^{\prime}}{\left\langle h_{\alpha_{j}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle}\right)=\left\langle h_{\alpha_{i}}^{\prime}, \frac{2 h_{\alpha_{j}}^{\prime}}{\left\langle h_{\alpha_{j}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle}\right\rangle=2 \frac{\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle}{\left\langle h_{\alpha_{j}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle}=A_{j i}
$$

Thus we obtain $\alpha_{i}=\sum_{j} A_{j i} \omega_{j}$.
In particular we note that all the fundamental roots are integral combinations of the fundamental weights, so lie in the weight lattice $X$. However, it is not true that the fundamental weights are, in general, integral combinations of the fundamental roots. We have

$$
\omega_{i}=\sum_{j}\left(A^{-1}\right)_{j i} \alpha_{j}
$$

For example, when $L$ has type $A_{1}$ we have

$$
\alpha_{1}=2 \omega_{1}, \quad \omega_{1}=\frac{1}{2} \alpha_{1} .
$$

When $L$ has type $A_{2}$ we have

$$
\begin{aligned}
& \alpha_{1}=2 \omega_{1}-\omega_{2} \\
& \alpha_{2}=-\omega_{1}+2 \omega_{2}
\end{aligned}
$$

and so

$$
\begin{aligned}
& \omega_{1}=\frac{2}{3} \alpha_{1}+\frac{1}{3} \alpha_{2} \\
& \omega_{2}=\frac{1}{3} \alpha_{1}+\frac{2}{3} \alpha_{2}
\end{aligned}
$$

We show that in general the coefficients expressing the $\omega_{i}$ in terms of the $\alpha_{j}$ are non-negative rational numbers.

Proposition 10.18 (i) $\left\langle\omega_{i}, \omega_{j}\right\rangle \geq 0$ for all $i, j$.
(ii) $\omega_{i}$ is a non-negative rational combination of $\alpha_{1}, \ldots, \alpha_{l}$.
(iii) The coefficients of the inverse $A^{-1}$ of the Cartan matrix are non-negative rational numbers.

Proof. We shall show that condition (i) implies the others, and prove (i) in a subsequent lemma. Since the coefficients of $A$ are integers the coefficients of $A^{-1}$ are rational numbers. We show they are all non-negative.

We know that $\omega_{i}\left(h_{j}\right)=\delta_{i j}$. This condition is equivalent to

$$
\left\langle\omega_{i}, \frac{2 \alpha_{j}}{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}\right\rangle=\delta_{i j}
$$

where $\langle$,$\rangle is now the Killing form on H^{*}$ as defined in Section 5.1. Thus we have

$$
\left\langle\omega_{i}, \alpha_{j}\right\rangle=\delta_{i j} \frac{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}{2}
$$

Now let $\omega_{i}=\sum_{j} c_{i j} \alpha_{j}$. Then we have

$$
\left\langle\omega_{i}, \omega_{j}\right\rangle=c_{i j}\left\langle\alpha_{j}, \omega_{j}\right\rangle=c_{i j} \frac{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}{2}
$$

Thus

$$
c_{i j}=2 \frac{\left\langle\omega_{i}, \omega_{j}\right\rangle}{\left\langle\alpha_{j}, \alpha_{j}\right\rangle} \geq 0
$$

since $\left\langle\omega_{i}, \omega_{j}\right\rangle \geq 0$ and $\left\langle\alpha_{j}, \alpha_{j}\right\rangle>0$.
We must now show that $\left\langle\omega_{i}, \omega_{j}\right\rangle \geq 0$. This will follow from the fact that

$$
\left\langle\omega_{i}, \frac{2 \alpha_{j}}{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}\right\rangle=\delta_{i j}
$$

and the fact that

$$
\left\langle\frac{2 \alpha_{i}}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}, \frac{2 \alpha_{j}}{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}\right\rangle \leq 0 \quad \text { if } i \neq j
$$

by Proposition 5.4. The following lemma on Euclidean spaces will give us what we need.

Lemma 10.19 Let $V$ be an n-dimensional Euclidean space with a basis $v_{1}, \ldots, v_{n}$ satisfying $\left\langle v_{i}, v_{j}\right\rangle \leq 0$ for all $i \neq j$. Let $w_{1}, \ldots, w_{n}$ be the dual basis of $V$ uniquely determined by the conditions $\left\langle v_{i}, w_{j}\right\rangle=\delta_{i j}$. Then $\left\langle w_{i}, w_{j}\right\rangle \geq 0$ for all $i, j$.

Proof. We use induction on $n$. If $n=1$ there is nothing to prove. So assume $n>1$ and let $U$ be the $(n-1)$-dimensional subspace of $V$ spanned by $v_{1}, \ldots, v_{n-1}$. Let $w_{1}^{\prime}, \ldots, w_{n-1}^{\prime}$ be the dual basis of $v_{1}, \ldots, v_{n-1}$ in $U$. Thus we have

$$
\left\langle v_{i}, w_{j}^{\prime}\right\rangle=\delta_{i j} \quad i, j=1, \ldots, n-1
$$

Let $U^{\perp}=\{v \in V ;\langle v, u\rangle=0$ for all $u \in U\}$. Then $\operatorname{dim} U^{\perp}=1$ and $U^{\perp}$ is the subspace of $V$ spanned by $w_{n}$. We see also that $w_{i}-w_{i}^{\prime} \in U^{\perp}$ for $i=1, \ldots, n-1$. Thus we have

$$
w_{i}=w_{i}^{\prime}+\lambda_{i} w_{n} \quad \text { for some } \lambda_{i} \in \mathbb{R}
$$

for $i=1, \ldots, n-1$. Taking the scalar product with $v_{n}$ we have

$$
0=\left\langle v_{n}, w_{i}^{\prime}\right\rangle+\lambda_{i}
$$

hence $\lambda_{i}=-\left\langle v_{n}, w_{i}^{\prime}\right\rangle$. We wish to determine the sign of $\lambda_{i}$. By induction we know $\left\langle w_{i}^{\prime}, w_{j}^{\prime}\right\rangle \geq 0$ for $i, j=1, \ldots, n-1$. This implies that $w_{i}^{\prime}$ is a non-negative combination of $v_{1}, \ldots, v_{n-1}$. Since

$$
\left\langle v_{n}, v_{1}\right\rangle \leq 0, \ldots\left\langle v_{n}, v_{n-1}\right\rangle \leq 0
$$

we see that $\left\langle v_{n}, w_{i}^{\prime}\right\rangle \leq 0$ and so $\lambda_{i} \geq 0$. Hence for $i, j=1, \ldots, n-1$

$$
\begin{aligned}
\left\langle w_{i}, w_{j}\right\rangle & =\left\langle w_{i}^{\prime}+\lambda_{i} w_{n}, w_{j}^{\prime}+\lambda_{j} w_{n}\right\rangle \\
& =\left\langle w_{i}^{\prime}, w_{j}^{\prime}\right\rangle+\lambda_{i} \lambda_{j}\left\langle w_{n}, w_{n}\right\rangle \geq 0
\end{aligned}
$$

since $\left\langle w_{i}^{\prime}, w_{j}^{\prime}\right\rangle \geq 0, \lambda_{i} \geq 0, \lambda_{j} \geq 0,\left\langle w_{n}, w_{n}\right\rangle>0$. It remains to show that $\left\langle w_{i}, w_{n}\right\rangle \geq 0$ for $i=1, \ldots, n-1$. We have

$$
\left\langle w_{i}, w_{n}\right\rangle=\left\langle w_{i}^{\prime}, w_{n}\right\rangle+\lambda_{i}\left\langle w_{n}, w_{n}\right\rangle \geq 0
$$

since $\left\langle w_{i}^{\prime}, w_{n}\right\rangle=0, \lambda_{i} \geq 0,\left\langle w_{n}, w_{n}\right\rangle>0$.
By applying this lemma in the case where $v_{i}=\frac{2 \alpha_{i}}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}, w_{i}=\omega_{i}$, we deduce that $\left\langle\omega_{i}, \omega_{j}\right\rangle \geq 0$ and so Proposition 10.18 is proved.

We now turn to the main theorem of the present section.

Theorem 10.20 Suppose $\lambda \in H^{*}$ is dominant and integral, that is $\lambda \in X^{+}$. Then the irreducible L-module $L(\lambda)$ is finite dimensional.

Proof. We have $L(\lambda)=M(\lambda) / J(\lambda)$. We know from Theorem 10.7 that $M(\lambda)$ is the direct sum of its weight spaces and from Theorem 10.9 that any submodule of $M(\lambda)$ is also the direct sum of its weight spaces. This applies in particular to $J(\lambda)$. It follows that $L(\lambda)=M(\lambda) / J(\lambda)$ is also the direct sum of its weight spaces. In fact the same proof as given in Theorem 10.9 shows that any $H$-submodule of $L(\lambda)$ is the direct sum of its weight spaces.

Let $v_{\lambda}$ be a highest weight vector of $L(\lambda)$. Thus $v_{\lambda} \in L(\lambda)_{\lambda}$ and $v_{\lambda} \neq 0$. We consider the sequence of elements

$$
v_{\lambda}, \quad f_{i} v_{\lambda}, \quad f_{i}^{2} v_{\lambda}, \quad \ldots
$$

We wish to show that terms in this sequence eventually become zero. In fact we show

$$
f_{i}^{k_{i}} v_{\lambda}=0 \quad \text { where } \quad k_{i}=\lambda\left(h_{i}\right)+1
$$

Let $m_{\lambda}$ be a highest weight vector of the Verma module $M(\lambda)$ such that $m_{\lambda}+J(\lambda)=v_{\lambda}$. We consider the submodule $\mathfrak{U}(L) f_{i}^{k_{i}} m_{\lambda}$ of $M(\lambda)$. As usual we choose elements $e_{i} \in L_{\alpha_{i}}, f_{i} \in L_{-\alpha_{i}}$ such that $\left[e_{i} f_{i}\right]=h_{i}$. We have

$$
\begin{aligned}
& e_{i} f_{i}=f_{i} e_{i}+h_{i} \\
& \begin{aligned}
e_{i} f_{i}^{2} & =f_{i} e_{i} f_{i}+h_{i} f_{i}=f_{i}^{2} e_{i}+2 f_{i} h_{i}-2 f_{i} \\
& =f_{i}^{2} e_{i}+2 f_{i}\left(h_{i}-1\right)
\end{aligned}
\end{aligned}
$$

and inductively we obtain

$$
e_{i} f_{i}^{n}=f_{i}^{n} e_{i}+n f_{i}^{n-1}\left(h_{i}-(n-1)\right) .
$$

Thus we have

$$
e_{i} f_{i}^{k_{i}} m_{\lambda}=f_{i}^{k_{i}} e_{i} m_{\lambda}+k_{i} f_{i}^{k_{i}-1}\left(h_{i}-\left(k_{i}-1\right)\right) m_{\lambda}=0
$$

since $e_{i} m_{\lambda}=0$ and $h_{i} m_{\lambda}=\lambda\left(h_{i}\right) m_{\lambda}=\left(k_{i}-1\right) m_{\lambda}$. Also if $j \neq i$ then

$$
e_{j} f_{i}^{k_{i}} m_{\lambda}=f_{i}^{k_{i}} e_{j} m_{\lambda}=0
$$

Thus $e_{j} f_{i}^{k_{i}} m_{\lambda}=0$ for all $j=1, \ldots, l$. It follows that $e_{\alpha} f_{i}^{k_{i}} m_{\lambda}=0$ for all $\alpha \in$ $\Phi^{+}$, since $e_{1}, \ldots, e_{l}$ generate $N$, by Proposition 7.7. We also know that

$$
h_{j} f_{i}^{k_{i}} m_{\lambda}=\left(\lambda-k_{i} \alpha_{i}\right)\left(h_{j}\right) f_{i}^{k_{i}} m_{\lambda}
$$

since $f_{i}^{k_{i}} m_{\lambda}$ is a weight vector with weight $\lambda-k_{i} \alpha_{i}$.

We now consider an arbitrary basis vector of $\mathfrak{H}(L)$ applied to $f_{i}^{k_{i}} m_{\lambda}$ :

$$
f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} e_{\beta_{1}}^{t_{1}} \cdots e_{\beta_{N}}^{t_{N}}\left(f_{i}^{k_{i}} m_{\lambda}\right)
$$

is zero unless all $t_{i}=0$, in which case it will be a scalar multiple of

$$
f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} f_{i}^{k_{i}} m_{\lambda}
$$

This shows that

$$
\mathfrak{U}(L) f_{i}^{k_{i}} m_{\lambda}=\mathfrak{U}\left(N^{-}\right) f_{i}^{k_{i}} m_{\lambda}
$$

Now $\mathfrak{U}\left(N^{-}\right) f_{i}^{k_{i}}$ is a proper subspace of $\mathfrak{U}\left(N^{-}\right)$since $k_{i}=\lambda\left(h_{i}\right)+1>0$. It follows from Theorem 10.6 that $\mathfrak{U}\left(N^{-}\right) f_{i}^{k_{i}} m_{\lambda}$ is a proper subspace of $M(\lambda)$. Hence $\mathfrak{H}(L) f_{i}^{k_{i}} m_{\lambda}$ is a proper submodule of $M(\lambda)$. It therefore lies in the unique maximal submodule $J(\lambda)$ of $M(\lambda)$. Hence $f_{i}^{k_{i}} m_{\lambda} \in J(\lambda)$ and this implies $f_{i}^{k_{i}} v_{\lambda}=0$.

Now let $K$ be the finite dimensional subspace of $L(\lambda)$ given by

$$
K=\mathbb{C} v_{\lambda}+\mathbb{C} f_{i} v_{\lambda}+\cdots+\mathbb{C} f_{i}^{k_{i}-1} v_{\lambda}
$$

We clearly have $H K \subset K$ since each $f_{i}^{n} v_{\lambda}$ is a weight vector. We have $f_{i} K \subset K$ since $f_{i}^{k_{i}} v_{\lambda}=0$. We also have $e_{i} K \subset K$ since

$$
\begin{aligned}
e_{i} f_{i}^{n} v_{\lambda} & =f_{i}^{n} e_{i} v_{\lambda}+n f_{i}^{n-1}\left(h_{i}-(n-1)\right) v_{\lambda} \\
& =n\left(\lambda\left(h_{i}\right)-(n-1)\right) f_{i}^{n-1} v_{\lambda}
\end{aligned}
$$

Thus $K$ is a submodule of $L(\lambda)$ for the subalgebra $\left\langle e_{i}, H, f_{i}\right\rangle$ of $L$ of dimension $l+2$. We shall consider non-zero finite dimensional $\left\langle e_{i}, H, f_{i}\right\rangle$-submodules of $L(\lambda)$. $K$ is such a submodule. If $U$ is any finite dimensional $\left\langle e_{i}, H, f_{i}\right\rangle$ submodule of $L(\lambda)$ we claim that $L U$ is also. For $L U$ is finite dimensional and we have, for $u \in U, z \in L, y \in\left\langle e_{i}, H, f_{i}\right\rangle$

$$
y(z u)=z(y u)+[y z] u \in L U
$$

since $y u \in U$ and $[y z] \in L$.
Let $V$ be the sum of all finite dimensional $\left\langle e_{i}, H, f_{i}\right\rangle$-submodules of $L(\lambda)$. Then $V \neq O$ since $V$ contains $K . V$ is an $L$-submodule of $L(\lambda)$, since if $U$ is a finite dimensional $\left\langle e_{i}, H, f_{i}\right\rangle$-submodule of $L(\lambda)$ so is $L U$. Since $L(\lambda)$ is an irreducible $L$-module we see that $V=L(\lambda)$. Thus $L(\lambda)$ is a sum of finite dimensional $\left\langle e_{i}, H, f_{i}\right\rangle$-submodules.

Now each such finite dimensional $\left\langle e_{i}, H, f_{i}\right\rangle$-submodule of $L(\lambda)$ is the direct sum of its weight spaces, as observed above. Thus we may choose a
basis for it consisting of weight vectors, that is vectors spanning 1-dimensional $H$-modules. Hence we can find a basis of $L(\lambda)$ consisting of weight vectors, each of which lies in some finite dimensional $\left\langle e_{i}, H, f_{i}\right\rangle$-submodule of $L(\lambda)$. This fact will give useful information about the set of weights of $L(\lambda)$.

Let $\Lambda$ be the set of all weights of $L(\lambda)$. Thus $\mu \in \Lambda$ if and only if $L(\lambda)_{\mu} \neq O$. Of course all weights of $L(\lambda)$ are weights of $M(\lambda)$ so have form

$$
\lambda-r_{1} \beta_{1}-\cdots-r_{N} \beta_{N}
$$

by Theorem 10.7. In particular $\Lambda \subset X$, since $\lambda \in X$ and each $\beta_{i} \in X$ by Theorem 10.7. Let $\mu$ be any element of $\Lambda$. Then there is a weight vector $v_{\mu} \in L(\lambda)$ for $\mu$ such that $v_{\mu}$ lies in a finite dimensional $\left\langle e_{i}, H, f_{i}\right\rangle$-submodule $U$ of $L(\lambda)$.

We consider the vectors

$$
\ldots, f_{i}^{2} v_{\mu}, f_{i} v_{\mu}, v_{\mu}, e_{i} v_{\mu}, e_{i}^{2} v_{\mu}, \ldots
$$

These vectors all lie in $U$ and have weights

$$
\ldots, \mu-2 \alpha_{i}, \mu-\alpha_{i}, \mu, \mu+\alpha_{i}, \mu+2 \alpha_{i}, \ldots
$$

Since $\operatorname{dim} U$ is finite $U$ has only finitely many weights so there exist $p, q \geq 0$ such that

$$
\begin{array}{ll}
f_{i}^{n} v_{\mu} \neq 0 \quad \text { for } 0 \leq n \leq p, & f_{i}^{p+1} v_{\mu}=0 \\
e_{i}^{n} v_{\mu} \neq 0 \quad \text { for } 0 \leq n \leq q, & e_{i}^{q+1} v_{\mu}=0
\end{array}
$$

Let $\quad V=\mathbb{C} f_{i}^{p} v_{\mu}+\cdots+\mathbb{C} f_{i} v_{\mu}+\mathbb{C} v_{\mu}+\mathbb{C} e_{i} v_{\mu}+\cdots+\mathbb{C} e_{i}^{q} v_{\mu}$. Then $V$ is a $\left\langle e_{i}, H, f_{i}\right\rangle$-submodule of $L(\lambda)$. This follows readily from the relations

$$
\begin{aligned}
e_{i} f_{i}^{n} & =f_{i}^{n} e_{i}+n f_{i}^{n-1}\left(h_{i}-(n-1)\right) \\
f_{i} e_{i}^{n} & =e_{i}^{n} f_{i}-n e_{i}^{n-1}\left(h_{i}+(n-1)\right)
\end{aligned}
$$

and the fact that $f_{i}^{p+1} v_{\mu}=0, e_{i}^{q+1} v_{\mu}=0$. We consider the trace of $h_{i}$ on $V$. On the one hand we have

$$
\begin{aligned}
\operatorname{trace}_{V} h_{i} & =\left(\mu\left(h_{i}\right)-p \alpha_{i}\left(h_{i}\right)\right)+\cdots+\mu\left(h_{i}\right)+\cdots+\left(\mu\left(h_{i}\right)+q \alpha_{i}\left(h_{i}\right)\right) \\
& =(p+q+1) \mu\left(h_{i}\right)+\left(\frac{q(q+1)}{2}-\frac{p(p+1)}{2}\right) \alpha_{i}\left(h_{i}\right) \\
& =(p+q+1) \mu\left(h_{i}\right)+(q-p)(p+q+1)
\end{aligned}
$$

since $\alpha_{i}\left(h_{i}\right)=2$. On the other hand we have

$$
\operatorname{trace}_{V} h_{i}=\operatorname{trace}_{V}\left[e_{i} f_{i}\right]=\operatorname{trace}_{V}\left(e_{i} f_{i}-f_{i} e_{i}\right)=0
$$

Hence $\mu\left(h_{i}\right)=p-q$.
We may make an important interpretation of this result in terms of the Weyl group $W$. We recall from Section 5.2 that $W$ is the group of linear transformations of $H_{\mathbb{R}}^{*}$ generated by the reflections $s_{\alpha}$ with respect to the roots $\alpha \in \Phi$. We write $s_{i}=s_{\alpha_{i}}$ and recall from Theorem 5.13 that $W$ is generated by $s_{1}, \ldots, s_{l}$. We have

$$
s_{i}(\mu)=\mu-2 \frac{\left\langle\alpha_{i}, \mu\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} \alpha_{i}=\mu-\mu\left(h_{i}\right) \alpha_{i}
$$

since

$$
\mu\left(h_{i}\right)=\mu\left(\frac{2 h_{\alpha_{i}}^{\prime}}{\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle}\right)=\left\langle\mu, \frac{2 \alpha_{i}}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}\right\rangle=2 \frac{\left\langle\alpha_{i}, \mu\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} .
$$

Choosing $\mu$ as above, where $\mu\left(h_{i}\right)=p-q$, we have

$$
s_{i}(\mu)=\mu-(p-q) \alpha_{i}=\mu+(q-p) \alpha_{i}
$$

Now $\mu+(q-p) \alpha_{i}$ is one of the weights in the list

$$
\mu-p \alpha_{i}, \ldots, \mu-\alpha_{i}, \mu, \mu+\alpha_{i}, \ldots, \mu+q \alpha_{i}
$$

of weights of $V$. Thus we have shown that if $\mu$ is any weight of $V$ then $s_{i}(\mu)$ is a weight of $L(\lambda)$ also. Since $s_{1}, \ldots, s_{l}$ generate $W$ it follows that for any $\mu \in \Lambda$ and any $w \in W$ we have $w(\mu) \in \Lambda$ also. Thus the set of weights $\Lambda$ of $L(\mu)$ is invariant under the Weyl group. We recall also from Proposition 5.8 that $W$ is finite.

We now claim that for each $\mu \in \Lambda$ there exists $w \in W$ such that $w(\mu) \in X^{+}$. To see this we consider the finite set of weights $\{w(\mu) ; w \in W\}$ and pick one maximal in the partial order $\succ$ on $H^{*}$. Let $\nu$ be such a weight. Then

$$
s_{i}(\nu)=\nu-\nu\left(h_{i}\right) \alpha_{i} .
$$

We know that $\nu \in X$ since $\Lambda \subset X$, hence $\nu\left(h_{i}\right) \in \mathbb{Z}$. If $\nu\left(h_{i}\right)<0$ we would have $s_{i}(\nu) \succ \nu$, a contradiction to the choice of $\nu$. Hence $\nu\left(h_{i}\right) \geq 0$. This holds for all $i=1, \ldots, l$ and so $\nu \in X^{+}$. Thus each weight in $\Lambda$ has a $W$-transform which lies in $X^{+}$.

We shall now concentrate on the set $\Lambda \cap X^{+}$. For any weight $\nu \in \Lambda \cap X^{+}$ we have $\nu \prec \lambda$. We express $\lambda$ and $\nu$ in terms of the fundamental roots $\alpha_{i}$. Since $\lambda, \nu \in X^{+}$these weights are non-negative integral combinations of the fundamental weights $\omega_{1}, \ldots, \omega_{l}$. By Proposition 10.18 they are therefore
non-negative rational combinations of the fundamental roots $\alpha_{1}, \ldots, \alpha_{l}$. Thus we have

$$
\begin{aligned}
& \lambda=\sum_{i=1}^{l} q_{i} \alpha_{i} \quad q_{i} \in \mathbb{Q} \quad q_{i} \geq 0 \\
& \nu=\sum_{i=1}^{l} q_{i}^{\prime} \alpha_{i} \quad q_{i}^{\prime} \in \mathbb{Q} \quad q_{i}^{\prime} \geq 0 .
\end{aligned}
$$

The condition $\lambda \succ \nu$ means simply that $q_{i}-q_{i}^{\prime}$ is a non-negative integer for each $i=1, \ldots, l$. Now given $q_{i}$ there are only finitely many $q_{i}^{\prime}$ such that $q_{i}^{\prime} \geq 0$ and $q_{i}-q_{i}^{\prime}$ is a non-negative integer. Thus given $\lambda \in X^{+}$there are only finitely many $\nu \in X^{+}$such that $\nu \prec \lambda$. Thus $\Lambda \cap X^{+}$is finite. Since every element of $\Lambda$ can be transformed by an element of $W$ into one of $\Lambda \cap X^{+}$and since $W$ is finite we see that $\Lambda$ is finite. Thus $L(\lambda)$ has only finitely many weights. However, each weight space $L(\lambda)_{\mu}$ of $L(\lambda)$ is finite dimensional, since

$$
\operatorname{dim} L(\lambda)_{\mu} \leq \operatorname{dim} M(\lambda)_{\mu}
$$

and $\operatorname{dim} M(\lambda)_{\mu}$ is finite by Theorem 10.7. Thus we have

$$
L(\lambda)=\bigoplus_{\mu} L(\lambda)_{\mu}
$$

with finitely many summands, each finite dimensional. Hence $L(\lambda)$ is finite dimensional.

We conclude by summarising the main ideas in this somewhat lengthy proof. In order to show that $L(\lambda)$ is finite dimensional it is sufficient to show that $L(\lambda)$ has only finitely many weights, since each weight space is known to be finite dimensional. This can be proved if the set of weights is known to be invariant under the Weyl group, since each weight will be $W$-equivalent to one in $X^{+}$, and there are only finitely many elements of $X^{+}$lower than $\lambda$ in the partial ordering. It is therefore necessary to show that, for any weight $\mu$ of $L(\lambda), s_{i}(\mu)$ is a weight also. This can be shown provided we know that any weight $\mu$ comes from a weight vector lying in a finite dimensional $\left\langle e_{i}, H, f_{i}\right\rangle$-submodule of $L(\lambda)$. We therefore have to show that $L(\lambda)$ is the sum of its finite dimensional $\left\langle e_{i}, H, f_{i}\right\rangle$-submodules. This comes from the irreducibility of $L(\lambda)$ provided $L(\lambda)$ has a non-zero finite dimensional $\left\langle e_{i}, H, f_{i}\right\rangle$-submodule. The existence of such a submodule $K$ is proved above.

We have now completed the determination of the finite dimensional irreducible $L$-modules where $L$ is a finite dimensional semisimple Lie algebra over $\mathbb{C}$.

Theorem 10.21 Let $L$ be a finite dimensional semisimple Lie algebra over $\mathbb{C}$. Then the finite dimensional irreducible L-modules are the modules $L(\lambda)$ for $\lambda \in X^{+}$. These modules are pairwise non-isomorphic.

Proof. The fact that any finite dimensional irreducible $L$-module is isomorphic to $L(\lambda)$ for some $\lambda$ is proved in Corollary 10.14. The fact that $\lambda$ must lie in $X^{+}$is proved in Proposition 10.15. The fact that $L(\lambda)$ is finite dimensional when $\lambda \in X^{+}$is proved in Theorem 10.20. The fact that the $L(\lambda)$ are pairwise non-isomorphic follows from the fact that $\lambda$ is the highest weight of $L(\lambda)$. Thus if $\lambda \neq \mu, L(\lambda)$ and $L(\mu)$ have different highest weights so cannot be isomorphic.

A property of $L(\lambda)$ which will be very useful subsequently is given by the following proposition.

Proposition 10.22 Let $\lambda \in X^{+}$and $w \in W$. Then

$$
\operatorname{dim} L(\lambda)_{\mu}=\operatorname{dim} L(\lambda)_{w(\mu)}
$$

Proof. Since $W$ is generated by the fundamental reflections $s_{1}, \ldots, s_{l}$ it is sufficient to show that

$$
\operatorname{dim} L(\lambda)_{\mu}=\operatorname{dim} L(\lambda)_{s_{i}(\mu)}
$$

We recall from Section 7.5 that there is an automorphism $\theta_{i}$ of $L$ such that $\theta_{i}(H)=H$ and $\theta_{i}(h)=s_{i}(h)$ for all $h \in H$. We define an $L$-module $\bar{L}(\lambda)$ which is the same space $L(\lambda)$ as before but with a different $L$-action. For $\bar{v} \in \bar{L}(\lambda)$ we have

$$
x \bar{v}=\overline{\theta_{i}(x) v}
$$

where $v$ is the corresponding element of $L(\lambda)$. It is clear that this action makes $\bar{L}(\lambda)$ into an $L$-module.

Now let $v \in L(\lambda)_{\mu}$. For $x \in H$ we have

$$
\begin{aligned}
x \bar{v} & =\overline{\theta_{i}(x) v}=\overline{s_{i}(x) v}=\overline{\mu\left(s_{i}(x)\right) v} \\
& =\left(s_{i}(\mu)\right)(x) \bar{v}
\end{aligned}
$$

Thus $\bar{v} \in \bar{L}(\lambda)_{s_{i}(\mu)}$. A similar argument shows that if $\bar{v} \in \bar{L}(\lambda)_{s_{i}(\mu)}$ then $v \in$ $L(\lambda)_{\mu}$. Hence

$$
\operatorname{dim} \bar{L}(\lambda)_{s_{i}(\mu)}=\operatorname{dim} L(\lambda)_{\mu}
$$

Now $\bar{L}(\lambda)$ is an irreducible $L$-module, since $L(\lambda)$ is irreducible. For if $\bar{M}$ were a submodule of $\bar{L}(\lambda)$ the corresponding subspace $M$ would be a submodule of $L(\lambda)$. Let $\Lambda$ be the set of weights of $L(\lambda)$. Then we have seen that $s_{i}(\Lambda)$ is the set of weights of $\bar{L}(\lambda)$. But we showed in the proof of Theorem 10.20 that $w(\Lambda)=\Lambda$ for all $w \in W$. Hence the set of weights of $\bar{L}(\lambda)$ is also $\Lambda$. In particular the highest weight of $\bar{L}(\lambda)$ is $\lambda$. Thus $\bar{L}(\lambda)$ is a finite dimensional irreducible $L$-module with highest weight $\lambda$. Hence $\bar{L}(\lambda)$ is isomorphic to $L(\lambda)$ by Theorem 10.21. Thus we have

$$
\operatorname{dim} L(\lambda)_{\mu}=\operatorname{dim} \bar{L}(\lambda)_{s_{i}(\mu)}=\operatorname{dim} L(\lambda)_{s_{i}(\mu)}
$$

Since each $w \in W$ is a product of elements $s_{i}$ we deduce that

$$
\operatorname{dim} L(\lambda)_{\mu}=\operatorname{dim} L(\lambda)_{w(\mu)}
$$

as required.

## 11

## Further properties of the universal enveloping algebra

### 11.1 Relations between the enveloping algebra and the symmetric algebra

Let $L$ be any finite dimensional Lie algebra over $\mathbb{C}$. Let $T$ be the tensor algebra of $L$. We recall that the enveloping algebra $\mathfrak{H}(L)$ is defined by

$$
\mathfrak{U l}(L)=T / J
$$

where $J$ is the 2 -sided ideal of $T$ generated by all elements of the form

$$
x \otimes y-y \otimes x-[x y]
$$

for $x, y \in L$. The symmetric algebra $S(L)$ is defined by

$$
S(L)=T / I
$$

where $I$ is the 2 -sided ideal of $T$ generated by all elements of the form

$$
x \otimes y-y \otimes x
$$

for $x, y \in L$.
$S(L)$ is isomorphic, as $\mathbb{C}$-algebra, to the polynomial ring $\mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$ where $n=\operatorname{dim} L$. We have

$$
S(L)=\bigoplus_{k} S^{k}(L)
$$

where $S^{k}(L)=\left(T^{k}+I\right) / I$.
$S^{k}(L)$ is the set of homogeneous elements of $S(L)$ of degree $k$. In particular we have an isomorphism

$$
L=T^{1} \rightarrow S^{1}(L)
$$

thus $L$ can be regarded as a subspace of $S(L)$.

If $x_{1}, \ldots, x_{n}$ are a basis of $L$ then the elements

$$
x_{1}^{r_{1}} \ldots x_{n}^{r_{n}} \quad r_{1}, \ldots, r_{n} \geq 0
$$

form a basis of $S(L)$.
We now explain how $S(L)$ can be regarded as a left $L$-module. In the first place $L$ is an $L$-module under the adjoint action. Then $T$ may be made into an $L$-module by means of the action

$$
y\left(x_{i_{1}} \otimes \cdots \otimes x_{i_{k}}\right)=\left[y x_{i_{1}}\right] \otimes x_{i_{2}} \otimes \cdots \otimes x_{i_{k}}+\cdots+x_{i_{1}} \otimes \cdots \otimes x_{i_{k-1}} \otimes\left[y x_{i_{k}}\right] .
$$

The ideal $I$ of $T$ is then a submodule, and so $S(L)=T / I$ can be given the structure of a left $L$-module. We have

$$
y\left(x_{i_{1}} \ldots x_{i_{k}}\right)=\left[y x_{i_{1}}\right] x_{i_{2}} \ldots x_{i_{k}}+\cdots+x_{i_{1}} \ldots x_{i_{k-1}}\left[y x_{i_{k}}\right]
$$

where $y \in L$ and the $x_{i_{\alpha}}$ are basis vectors of $L$. We note that each $S^{k}(L)$ is an $L$-submodule of $S(L)$.

Similarly $\mathfrak{U}(L)=T / J$ can be made into a left $L$-module. For the ideal $J$ of $T$ is also a submodule since, for $a, b \in L$, we have

$$
\begin{aligned}
y(a \otimes b-b \otimes a-[a b])= & {[y a] \otimes b+a \otimes[y b]-[y b] \otimes a-b \otimes[y a]-[y[a b]] } \\
= & {[y a] \otimes b-b \otimes[y a]-[[y a] b]+a \otimes[y b] } \\
& -[y b] \otimes a-a[y b]
\end{aligned}
$$

since

$$
[y[a b]]=[[y a] b]+[a[y b]] .
$$

We shall find it useful to compare the enveloping algebra $\mathfrak{H}(L)$ with the symmetric algebra $S(L)$. We first compare their $\mathbb{C}$-algebra structures. Of course they need not be isomorphic as $\mathbb{C}$-algebras since $S(L)$ is commutative whereas $\mathfrak{H}(L)$ is in general non-commutative. However, there is a relation between these two algebras: it is the relation between a filtered algebra and the corresponding graded algebra.

A filtered algebra is an associative algebra $A$ with a chain of subspaces

$$
A_{0} \subset A_{1} \subset A_{2} \subset \cdots
$$

such that $\cup_{i} A_{i}=A$ and $A_{i} A_{j} \subset A_{i+j}$.
A graded algebra is an associative algebra $A$ with a decomposition

$$
A=A_{0} \oplus A_{1} \oplus A_{2} \oplus \cdots
$$

into a direct sum of subspaces such that $A_{i} A_{j} \subset A_{i+j} \quad$ for all $i, j$.

Given any filtered algebra we may obtain a corresponding graded algebra as follows. Let $A=\bigcup_{i} A_{i}$ be a filtered algebra. We define vector spaces $B_{0}$, $B_{1}, B_{2}, \ldots$ by

$$
B_{0}=A_{0}, \quad B_{1}=A_{1} / A_{0}, \quad B_{2}=A_{2} / A_{1}, \quad \ldots
$$

and define the vector space $B$ by

$$
B=B_{0} \oplus B_{1} \oplus B_{2} \oplus \cdots
$$

We define a multiplication on $B$ to make it into a graded algebra. It is sufficient to define $x y$ when $x \in B_{i}, y \in B_{j}$ and to extend this multiplication by linearity. Thus let $x \in A_{i} / A_{i-1}, y \in A_{j} / A_{j-1}$. Let $x=A_{i-1}+a_{i}, y=A_{j-1}+a_{j}$. Then, for any pair of elements $u \in A_{i-1}, v \in A_{j-1}$ we have

$$
\left(u+a_{i}\right)\left(v+a_{j}\right)=u v+u a_{j}+a_{i} v+a_{i} a_{j} \in A_{i+j-1}+a_{i} a_{j}
$$

Thus the coset in $A_{i+j} / A_{i+j-1}$ containing the product of any element in $x$ with any element in $y$ is the same. Thus we may without ambiguity define $x y \in B_{i+j}$ by

$$
x y=A_{i+j-1}+a_{i} a_{j}
$$

It is readily checked that this multiplication when extended by linearity makes $B$ into a graded algebra. $B$ is called the associated graded algebra of the filtered algebra $A$.

We may regard $\mathfrak{l}(L)$ as a filtered algebra as follows. Let $\mathfrak{U}_{i}(L)$ be the subspace of $\mathfrak{H}(L)$ generated by all products $a_{1} a_{2} \ldots a_{j}$ for $j \leq i$, where $a_{k} \in L$. We also define $\mathfrak{U}_{0}(L)=\mathbb{C} 1$. Then we have

$$
\bigcup_{i} \mathfrak{l}_{i}(L)=\mathfrak{l}(L)
$$

and

$$
\mathfrak{U}_{0}(L) \subset \mathfrak{H}_{1}(L) \subset \mathfrak{H}_{2}(L) \subset \cdots
$$

Moreover $\mathfrak{U}_{i}(L) \mathfrak{U}_{j}(L) \subset \mathfrak{U}_{i+j}(L)$. Thus $\mathfrak{H}(L)$ is a filtered algebra. We consider its associated graded algebra.

Proposition 11.1 The associated graded algebra of the filtered algebra $\mathfrak{H}(L)$ is isomorphic to $S(L)$.

Proof. Let $B=B_{0} \oplus B_{1} \oplus B_{2} \oplus \cdots$ be the associated graded algebra of $\mathfrak{H}(L)$. We first observe that $B$ is a commutative algebra. $B$ is generated as an algebra by 1 and $B_{1}$, and $B_{1}=\mathfrak{U}_{1}(L) / \mathfrak{H}_{0}(L)$. The natural map

$$
L \rightarrow \mathfrak{U}_{1}(L) / \mathfrak{U}_{0}(L)
$$

is an isomorphism of vector spaces. For elements $x, y \in L$ we have

$$
x y-y x=[x y] \quad \text { in } \quad \mathfrak{U}(L)
$$

Thus

$$
\left(\mathfrak{U}_{0}(L)+x\right)\left(\mathfrak{U}_{0}(L)+y\right) \equiv\left(\mathfrak{U}_{0}(L)+y\right)\left(\mathfrak{U}_{0}(L)+x\right) \quad \bmod \mathfrak{U}_{1}(L)
$$

Hence any two elements of $B_{1}=\mathfrak{U}_{1}(L) / \mathfrak{U}_{0}(L)$ commute in $B$, where their product lies in $\mathfrak{U}_{2}(L) / \mathfrak{U}_{1}(L)$. It follows that $B$ is a commutative algebra.

We now compare $B$ with the symmetric algebra $S(L)$. Let $x_{1}, \ldots, x_{n}$ be a basis of $L$. Then it follows from the PBW basis theorem that the elements

$$
x_{1}^{r_{1}} \ldots x_{n}^{r_{n}} \quad r_{1}+\cdots+r_{n} \leq i
$$

form a basis of $\mathfrak{H}_{i}(L)$. Moreover the elements $\mathfrak{1}_{i-1}(L)+x_{1}^{r_{1}} \ldots x_{n}^{r_{n}}$ with $r_{1}+$ $\cdots+r_{n}=i$ form a basis for $\mathfrak{U}_{i}(L) / \mathfrak{U}_{i-1}(L)=B_{i}$. Now we have

$$
\left(\mathfrak{U}_{i-1}(L)+x_{1}^{r_{1}} \ldots x_{n}^{r_{n}}\right)\left(\mathfrak{U}_{j-1}(L)+x_{1}^{s_{1}} \ldots x_{n}^{s_{n}}\right)=\mathfrak{U}_{i+j-1}(L)+x_{1}^{r_{1}} \ldots x_{n}^{r_{n}} x_{1}^{s_{1}} \ldots x_{n}^{s_{n}} .
$$

This is equal to

$$
\mathfrak{U}_{i+j-1}(L)+x_{1}^{r_{1}+s_{1}} \ldots x_{n}^{r_{n}+s_{n}}
$$

since multiplication in $B$ is commutative. This shows that the linear map $S(L) \rightarrow B$ defined by

$$
x_{1}^{r_{1}} \ldots x_{n}^{r_{n}} \rightarrow \mathfrak{U}_{i-1}(L)+x_{1}^{r_{1}} \ldots x_{n}^{r_{n}} \quad \sum r_{k}=i
$$

extends to an isomorphism of algebras. Thus the associated graded algebra of $\mathfrak{H}(L)$ is isomorphic to $S(L)$.

We now wish to compare the enveloping algebra $\mathfrak{H}(L)$ and the symmetric algebra $S(L)$ as left $L$-modules. We shall show that they are isomorphic as $L$-modules. In order to do so we shall first find a complement to $\mathfrak{U}_{i-1}(L)$ in $\mathfrak{U}_{i}(L)$.

We have $T^{i}=L \otimes \cdots \otimes L$ ( $i$ factors). The symmetric group $S_{i}$ operates on $T^{i}$ by

$$
\sigma\left(y_{1} \otimes \cdots \otimes y_{i}\right)=y_{\sigma^{-1}(1)} \otimes \cdots \otimes y_{\sigma^{-1}(i)}
$$

and extending by linearity. A tensor in $T^{i}$ is called symmetric if it is fixed by all $\sigma \in S_{i}$. The natural map $T \rightarrow \mathfrak{U}(L)$ induces a map $T^{i} \rightarrow \mathfrak{U}_{i}(L)$. Let $\mathfrak{U}^{i}(L)$ be the image under this map of the space of symmetric tensors in $T^{i}$.

Proposition 11.2 (i) $\mathfrak{U}_{i}(L)=\mathfrak{U}_{i-1}(L) \oplus \mathfrak{U}^{i}(L)$.
(ii) These spaces are all L-submodules of $\mathfrak{H}(L)$.

Proof. We first show that

$$
\mathfrak{U}_{i}(L)=\mathfrak{U}_{i-1}(L)+\mathfrak{U}^{i}(L)
$$

Let $x_{1}^{r_{1}} \ldots x_{n}^{r_{n}}$ be a basis element of $\mathfrak{U}_{i}(L)$ with $r_{1}+\cdots+r_{n}=i$. For each $\sigma \in S_{i}$ we define $\sigma\left(x_{1}^{r_{1}} \ldots x_{n}^{r_{n}}\right)$ to be the element obtained from $x_{1}^{r_{1}} \ldots x_{n}^{r_{n}}$ by permuting the factors by the permutation $\sigma$. Since multiplication in the graded algebra of $\mathfrak{U}(L)$ is commutative we have

$$
x_{1}^{r_{1}} \ldots x_{n}^{r_{n}}=\frac{1}{i!} \sum_{\sigma \in S_{i}} \sigma\left(x_{1}^{r_{1}} \ldots x_{n}^{r_{n}}\right)+u
$$

where $u \in \mathfrak{U}_{i-1}(L)$. Since the sum lies in $\mathfrak{H}^{i}(L)$ we have

$$
\mathfrak{U}_{i}(L)=\mathfrak{U}_{i-1}(L)+\mathfrak{U}^{i}(L)
$$

We next show that $\mathfrak{U}_{i-1}(L) \cap \mathfrak{U}^{i}(L)=O$. Any element of $\mathfrak{U}^{i}(L)$ has the form

$$
\sum_{\substack{r_{1}, \ldots, r_{n} \\ r_{1}+\ldots+r_{n}=i}} \lambda_{r_{1}, \ldots, r_{n}} \sum_{\sigma \in S_{i}} \sigma\left(x_{1}^{r_{1}} \ldots x_{n}^{r_{n}}\right) .
$$

We express this element as a linear combination of basis elements of $\mathfrak{H}(L)$. We obtain

$$
\sum_{\substack{r_{1}, \ldots, r_{n} \\ r_{1}+\cdots+r_{n}=i}} \lambda_{r_{1}, \ldots, r_{n}} \sum_{\sigma \in S_{i}} \sigma\left(x_{1}^{r_{1}} \ldots x_{n}^{r_{n}}\right)=i!\sum_{\substack{r_{1}, \ldots, r_{n} \\ r_{1}+\cdots+r_{n}=i}} \lambda_{r_{1}, \ldots, r_{n}} x_{1}^{r_{1}} \ldots x_{n}^{r_{n}}+u
$$

where $u \in \mathfrak{U}_{i-1}(L)$, since multiplication in the graded algebra of $\mathfrak{H}(L)$ is commutative. This element can only lie in $\mathfrak{U}_{i-1}(L)$ if each $\lambda_{r_{1}}, \ldots,{ }_{r_{n}}$ is 0 . Thus $\mathfrak{U}_{i-1}(L) \cap \mathfrak{U}^{i}(L)=O$. Hence we have

$$
\mathfrak{U}_{i}(L)=\mathfrak{U}_{i-1}(L) \oplus \mathfrak{U}^{i}(L)
$$

Finally these subspaces are all $L$-submodules. The subspaces $\mathfrak{U}_{i}(L)$ and $\mathfrak{U}_{i-1}(L)$ are evidently submodules by the definition of the $L$-action. $\mathfrak{U}^{i}(L)$ is an $L$-submodule since the $L$-action commutes with the $S_{i}$-action on $T^{i}$.

Let $T_{\text {sym }}^{i}$ be the subspace of symmetric tensors in $T^{i}$.

Proposition 11.3 There is a commutative diagram of vector space isomorphisms
where $\alpha$ is induced by the map $T(L) \rightarrow \mathfrak{H}(L), \beta$ is induced by $T(L) \rightarrow S(L)$, $\gamma$ is induced by $\mathfrak{U}_{i}(L) \rightarrow \mathfrak{U}_{i}(L) / \mathfrak{U}_{i-1}(L)$ and $\delta$ is the map of Proposition 11.1.

## Example

Proof. It is sufficient to show that $\gamma \alpha(t)=\delta \beta(t)$ where $t=\sum_{\sigma \in S_{i}} y_{\sigma^{-1}(1)} \otimes$ $\cdots \otimes y_{\sigma^{-1}(i)}$ and $y_{k} \in L$. We have

$$
\begin{aligned}
\alpha(t) & =\sum_{\sigma} y_{\sigma^{-1}(1)} \ldots y_{\sigma^{-1}(i)} \\
\gamma \alpha(t) & =\mathfrak{l}_{i-1}(L)+\sum_{\sigma} y_{\sigma^{-1}(1)} \ldots y_{\sigma^{-1}(i)} \\
\beta(t) & =\sum_{\sigma} y_{\sigma^{-1}(1)} \ldots y_{\sigma^{-1}(i)} \\
\delta \beta(t) & =\mathfrak{l}_{i-1}(L)+\sum_{\sigma} y_{\sigma^{-1}(1)} \ldots y_{\sigma^{-1}(i)}
\end{aligned}
$$

since the difference between $y_{\sigma^{-1}(1)} \ldots y_{\sigma^{-1}(i)}$ and the corresponding element in canonical form lies in $\mathfrak{U}_{i-1}(L)$. Hence $\gamma \alpha(t)=\delta \beta(t)$.

We now define $\theta: S^{i}(L) \rightarrow \mathfrak{U}^{i}(L)$ by $\theta=\gamma^{-1} \delta$, and extend this map by linearity to give $\theta: S(L) \rightarrow \mathfrak{H}(L)$. $\theta$ is called the operation of symmetrisation. We have

$$
\theta\left(y_{1} y_{2} \ldots y_{i}\right)=\frac{1}{i!} \sum_{\sigma \in S_{i}} y_{\sigma^{-1}(1)} \ldots y_{\sigma^{-1}(i)}
$$

Proposition $11.4 \theta: S(L) \rightarrow \mathfrak{H}(L)$ is an isomorphism of L-modules.

Proof. We know that $\theta$ is an isomorphism of vector spaces and so must show that

$$
x \cdot \theta(P)=\theta(x \cdot P) \quad \text { for all } x \in L, P \in S(L)
$$

A derivation of an associative algebra $A$ is a linear map $D: A \rightarrow A$ such that

$$
D(a b)=D(a) b+a D(b)
$$

for all $a, b \in A$. It follows from the definition of the $L$-action that the maps

$$
\begin{array}{rlrl}
S(L) & \rightarrow S(L) & \mathfrak{H}(L) & \rightarrow \mathfrak{H}(L) \\
P & \rightarrow x \cdot P & u & \rightarrow x \cdot u
\end{array}
$$

for $x \in L$ are derivations. Now $L$ may be identified with a subspace of $S(L)$ and the map $P \rightarrow x \cdot P$ when restricted to $L$ is ad $x$. Similarly $L$ may be identified with a subspace of $\mathfrak{H}(L)$ and the map $u \rightarrow x \cdot u$ when restricted to $L$ is again ad $x$. Now $S(L)$ is generated as an algebra by $L$ and 1 . We have $D(1)=0$ for any derivation of $S(L)$. Thus there is a unique derivation of $S(L)$ extending ad $x$ on $L$. Similarly $u \rightarrow x \cdot u$ is the unique derivation of $\mathfrak{H}(L)$ extending ad $x$ on $L$.

Let $D: \mathfrak{U}(L) \rightarrow \mathfrak{U}(L)$ be this derivation. $D$ transforms $\mathfrak{U}^{i}(L)$ into $\mathfrak{U}^{i}(L)$ for each $i$. Using the isomorphism $\gamma$ of Proposition 11.3, $D$ determines a map

$$
\bigoplus_{i} \frac{\mathfrak{U}_{i}(L)}{\mathfrak{U}_{i-1}(L)} \rightarrow \bigoplus_{i} \frac{\mathfrak{U}_{i}(L)}{\mathfrak{U}_{i-1}(L)}
$$

which is still a derivation. Using the isomorphism $\delta$ of Proposition 11.3 we obtain a map $S(L) \rightarrow S(L)$ that is still a derivation and which acts as ad $x$ on $L$. Thus it is the map $P \rightarrow x \cdot P$. Hence for $P \in S(L)$ we have

$$
\theta^{-1}(x \cdot \theta(P))=x \cdot P
$$

Thus $x \cdot \theta(P)=\theta(x \cdot P)$ as required.
Note The $L$-action on $\mathfrak{H}(L)$ considered here may be described simply by

$$
x \cdot u=x u-u x \quad x \in L, u \in \mathfrak{H}(L)
$$

For this is a derivation of $\mathfrak{H}(L)$ which extends ad $x: L \rightarrow L$.

### 11.2 Invariant polynomial functions

Let $G=\operatorname{Inn} L$ be the group of inner automorphisms of the Lie algebra $L$. We recall from Section 3.2 that $G$ is generated by automorphisms of the form $\exp$ ad $x$ for elements $x \in L$ such that ad $x$ is nilpotent. We define an action of $G$ on $L^{*}$ by

$$
(g f) x=f\left(g^{-1} x\right) \quad g \in G, f \in L^{*}, x \in L
$$

The tensor algebra

$$
T\left(L^{*}\right)=\bigoplus_{k \geq 0}\left(L^{*} \underset{k \text { factors }}{\left.\otimes \cdots \otimes L^{*}\right)}\right.
$$

may then be made into a $G$-module satisfying

$$
g\left(f_{1} \otimes \cdots \otimes f_{k}\right)=g f_{1} \otimes \cdots \otimes g f_{k}
$$

for $g \in G, f_{i} \in L^{*}$. Let $I$ be the 2-sided ideal of $T\left(L^{*}\right)$ generated by all elements of form

$$
f \otimes g-g \otimes f
$$

for $f, g \in L^{*}$. Then $I$ is a $G$-submodule of $T\left(L^{*}\right)$. Let

$$
S\left(L^{*}\right)=T\left(L^{*}\right) / I .
$$

Then $S\left(L^{*}\right)$ may also be made into a $G$-module. $S\left(L^{*}\right)$ is the symmetric algebra on $L^{*}$. The algebra $S\left(L^{*}\right)$ may be identified with the algebra of polynomial functions on $L$. The element $I+f_{1} \otimes \cdots \otimes f_{k}$ of $S\left(L^{*}\right)$ gives rise to the polynomial function $f_{1} f_{2} \ldots f_{k}$ on $L$. We define $P(L)=S\left(L^{*}\right)$ and $P^{m}(L)=S^{m}\left(L^{*}\right)$. This is the image of $T^{m}\left(L^{*}\right)$ under the natural homomorphism $T\left(L^{*}\right) \rightarrow S\left(L^{*}\right) . P^{k}(L)$ is the space of homogeneous polynomial functions of degree $k$ on $L$. In particular $P^{1}(L)$ may be identified with $L^{*}$. Each subspace $P^{k}(L)$ is clearly a $G$-submodule of $P(L)$.

We now prove some lemmas which will help in understanding the action of $G$ on $P(L)$.

Lemma 11.5 The linear map $\alpha: T^{m}\left(L^{*}\right) \rightarrow\left(T^{m}(L)\right)^{*}$ uniquely determined by

$$
\left(\alpha\left(f_{1} \otimes \cdots \otimes f_{m}\right)\right)\left(x_{1} \otimes \cdots \otimes x_{m}\right)=f_{1}\left(x_{1}\right) f_{2}\left(x_{2}\right) \ldots f_{m}\left(x_{m}\right)
$$

is an isomorphism of $G$-modules. Here $x_{1}, \ldots, x_{m}$ lie in $L$ and $f_{1}, \ldots, f_{m}$ in $L^{*}$.

Proof. The linear map $\alpha$ is clearly injective. Since $T^{m}\left(L^{*}\right)$ and $\left(T^{m}(L)\right)^{*}$ have the same dimension, $\alpha$ must also be surjective. Thus $\alpha$ is an isomorphism of vector spaces. We must also show that

$$
\alpha\left(\gamma \cdot f_{1} \otimes \cdots \otimes f_{m}\right)=\gamma\left(\alpha\left(f_{1} \otimes \cdots \otimes f_{m}\right)\right)
$$

for all $\gamma \in G$. Now we have

$$
\gamma\left(f_{1} \otimes \cdots \otimes f_{m}\right)=\gamma f_{1} \otimes \cdots \otimes \gamma f_{m}
$$

Thus

$$
\begin{aligned}
\alpha\left(\gamma \cdot f_{1} \otimes \cdots \otimes f_{m}\right)\left(x_{1} \otimes \cdots \otimes x_{m}\right) & =\left(\gamma f_{1}\right)\left(x_{1}\right) \ldots\left(\gamma f_{m}\right)\left(x_{m}\right) \\
& =f_{1}\left(\gamma^{-1} x_{1}\right) \ldots f_{m}\left(\gamma^{-1} x_{m}\right)
\end{aligned}
$$

On the other hand

$$
\begin{aligned}
\gamma\left(\alpha\left(f_{1} \otimes \cdots \otimes f_{m}\right)\right)\left(x_{1} \otimes \cdots \otimes x_{m}\right) & =\alpha\left(f_{1} \otimes \cdots \otimes f_{m}\right)\left(\gamma^{-1} x_{1} \otimes \cdots \otimes \gamma^{-1} x_{m}\right) \\
& =f_{1}\left(\gamma^{-1} x_{1}\right) \ldots f_{m}\left(\gamma^{-1} x_{m}\right)
\end{aligned}
$$

This gives the required equality.

Lemma 11.6 Consider the maps

$$
T^{m}(L)^{*} \underset{\alpha^{-1}}{\longrightarrow} T^{m}\left(L^{*}\right) \underset{\beta}{\longrightarrow} S^{m}\left(L^{*}\right)
$$

and let $\theta:\left(T^{m}(L)\right)^{*} \longrightarrow S^{m}\left(L^{*}\right)$ be given by $\theta=\beta \alpha^{-1}$. Thus $\theta$ is a homomorphism of $G$-modules. Then we have $(\theta f) x=f(x \otimes \cdots \otimes x)$ with $m$ factors, for $x \in L$.

Proof. It is sufficient to prove this when $f$ has the form

$$
f\left(x_{1} \otimes \cdots \otimes x_{m}\right)=f_{1}\left(x_{1}\right) \ldots f_{m}\left(x_{m}\right)
$$

that is when $\alpha^{-1} f=f_{1} \otimes \cdots \otimes f_{m}$. In this case we have

$$
(\theta f) x=f_{1}(x) f_{2}(x) \ldots f_{m}(x)=f(x \otimes \cdots \otimes x)
$$

An element $f \in\left(T^{m}(L)\right)^{*}$ is called symmetric if

$$
f\left(x_{1} \otimes \cdots \otimes x_{m}\right)=f\left(x_{\sigma(1)} \otimes \cdots \otimes x_{\sigma(m)}\right)
$$

for all $x_{1}, \ldots, x_{m} \in L$ and all $\sigma \in S_{m}$. The set of symmetric elements of $\left(T^{m}(L)\right)^{*}$ will be denoted by $\left(T^{m}(L)\right)_{\text {sym }}^{*}$.

An element of $T^{m}\left(L^{*}\right)$ is called symmetric if it is invariant under the linear maps which transform $f_{1} \otimes \cdots \otimes f_{m}$ to $f_{\sigma(1)} \otimes \cdots \otimes f_{\sigma(m)}$ for all $\sigma \in S_{m}$.

The set of symmetric elements of $T^{m}\left(L^{*}\right)$ will be denoted by $T^{m}\left(L^{*}\right)_{\text {sym }}$.
Lemma 11.7 The subspaces $T^{m}(L)_{\text {sym }}^{*}$ and $T^{m}\left(L^{*}\right)_{\mathrm{sym}}$ are $G$-submodules. Moreover the maps $\alpha^{-1}, \beta$ give isomorphisms

$$
T^{m}(L)_{\mathrm{sym}}^{*} \underset{\alpha-1}{\longrightarrow} T^{m}\left(L^{*}\right)_{\mathrm{sym}} \xrightarrow[\beta]{\longrightarrow} S^{m}\left(L^{*}\right)
$$

Proof. The subspaces are $G$-submodules since the $G$-action commutes with the $S_{m}$-action on $T^{m}(L)^{*}$ and $T^{m}\left(L^{*}\right)$. The map $\alpha$ transforms $T^{m}\left(L^{*}\right)_{\text {sym }}$ into $T^{m}(L)_{\text {sym }}^{*}$ and, since $\alpha$ is an isomorphism and these two spaces have the same dimension, we have

$$
\alpha\left(T^{m}\left(L^{*}\right)_{\mathrm{sym}}\right)=T^{m}(L)_{\mathrm{sym}}^{*}
$$

Again the spaces $T^{m}\left(L^{*}\right)_{\text {sym }}$ and $S^{m}\left(L^{*}\right)$ have the same dimension and the map

$$
\beta: T^{m}\left(L^{*}\right)_{\mathrm{sym}} \longrightarrow S^{m}\left(L^{*}\right)
$$

is surjective, since it transforms

$$
\frac{1}{m!} \sum_{\sigma \in S_{m}} f_{\sigma(1)} \otimes \cdots \otimes f_{\sigma(m)}
$$

into $f_{1} f_{2} \ldots f_{m}$. Thus this map is also an isomorphism.

The $G$-module isomorphism

$$
P^{m}(L) \underset{\beta^{-1} \alpha}{\longrightarrow} T^{m}(L)_{\mathrm{sym}}^{*}
$$

is useful in determining the $G$-action on $P^{m}(L)$, since it is often easier to calculate the action on the linear functions in $T^{m}(L)_{\text {sym }}^{*}$ than on the polynomial functions in $P^{m}(L)$.

We shall now assume that the Lie algebra $L$ is semisimple. The group $G$ of inner automorphisms is called the adjoint group of $L$. A polynomial function $P \in P(L)$ is called invariant if $\gamma(P)=P$ for all $\gamma \in G$. The set of invariant polynomial functions on $L$ is denoted by $P(L)^{G}$. This is clearly a subalgebra of $P(L)$. We shall investigate the algebra of invariant polynomial functions on $L$ by relating it to the algebra of polynomial functions on a Cartan subalgebra of $L$ invariant under the Weyl group.

Let $H$ be a Cartan subalgebra of $L$ and $P(H)=S\left(H^{*}\right)$ be the algebra of polynomial functions on $H$. Let $W$ be the Weyl group of $L$. Then we know that both $H$ and $H^{*}$ are $W$-modules. (We recall from Section 5.2 that an action of $W$ was defined on the real subspace $H_{\mathbb{R}}^{*}$ of $H^{*}$, and this gives rise to a $W$-action on $H^{*}$ by linearity.) The $W$-actions on $H$ and $H^{*}$ are related by

$$
(w f) h=f\left(w^{-1} h\right) \quad w \in W, f \in H^{*}, h \in H
$$

There is then a $W$-action on $T\left(H^{*}\right)$ satisfying

$$
w\left(f_{1} \otimes \cdots \otimes f_{m}\right)=w f_{1} \otimes \cdots \otimes w f_{m}
$$

This in turn induces a $W$-action on $S\left(H^{*}\right)=T\left(H^{*}\right) / I$ since $I$ is a $W$-submodule of $T\left(H^{*}\right)$. Thus $P(H)=S\left(H^{*}\right)$ may be regarded as a $W$-module.

A polynomial function $P \in P(H)$ is called $W$-invariant if $w(P)=P$ for all $w \in W$. The set of $W$-invariant polynomial functions on $H$ will be denoted by $P(H)^{W}$. This is a subalgebra of $P(H)$.

Now we have an algebra homomorphism

$$
\psi: P(L) \rightarrow P(H)
$$

given by restriction from $L$ to $H$. We consider the image of $P(L)^{G}$ under this restriction map. We show first that this image lies in the subalgebra $P(H)^{W}$.

Proposition $11.8 \psi\left(P(L)^{G}\right) \subset P(H)^{W}$.

Proof. We use the element $\theta_{i} \in G$ given by

$$
\theta_{i}=\exp \text { ad } e_{i} \cdot \exp \operatorname{ad}\left(-f_{i}\right) \cdot \exp \text { ad } e_{i}
$$

We recall from Proposition 7.18 that $\theta_{i}(H)=H$ and that $\theta_{i}(h)=s_{i}(h)$ for all $h \in H$, where $s_{i} \in W$ is a fundamental reflection. Thus $\theta_{i}$ acts on $H$ in the same way as $s_{i}$. It follows that $\theta_{i}$ and $s_{i}$ also act in the same way on $H^{*}$, and on $S\left(H^{*}\right)=P(H)$.

Let $P \in P(L)^{G}$. Then $\theta_{i}(P)=P$. We have $\psi(P) \in P(H)$ and so $s_{i}(\psi(P))=$ $\psi(P)$. However, the Weyl group $W$ is generated by its fundamental reflections $s_{1}, \ldots, s_{l}$, by Theorem 5.13. Thus we have

$$
w(\psi(P))=\psi(P) \quad \text { for all } w \in W
$$

and so $\psi(P) \in P(H)^{W}$.
Proposition 11.9 The map $\psi: P(L)^{G} \rightarrow P(H)^{W}$ is injective.

Proof. Let $R$ be the set of regular elements of $L$. We recall from the proof of Proposition 3.12 that there is a polynomial function $F \in P(L)$ such that $x \in R$ if and only if $F(x) \neq 0$. We also recall from Theorem 3.2 that every regular element lies in some Cartan subalgebra and from Theorem 3.13 that any two Cartan subalgebras are conjugate. Thus given any regular element $x \in R$ there exists $\gamma \in G$ such that $\gamma(x) \in H$.

Now suppose $P \in P(L)^{G}$ satisfies $\psi(P)=O$. Let $x$ be a regular element and let $\gamma \in G$ be such that $\gamma(x) \in H$. Then

$$
\psi(P)(\gamma(x))=0
$$

that is $P(\gamma(x))=0$. Hence $\left(\gamma^{-1} P\right)(x)=0$. Since $P \in P(L)^{G}$ we have $\gamma^{-1} P=P$ and we may deduce that

$$
P(x)=0
$$

Thus $P$ annihilates all regular elements of $L$. Hence $P(x)=0$ whenever $F(x) \neq 0$. By the principle of irrelevance of algebraic inequalities we have

$$
P(x)=0 \quad \text { for all } x \in L
$$

that is $P=O$.
Finally we show that the map $\psi$ is also surjective.

Theorem $11.10 \psi: P(L)^{G} \rightarrow P(H)^{W}$ is surjective, and is therefore an isomorphism of algebras.

Proof. We make use of ideas from the representation theory of $L$. Let $\lambda \in H^{*}$ be a dominant integral weight and $L(\lambda)$ be the finite dimensional irreducible $L$-module with highest weight $\lambda$. We can choose a basis of $L(\lambda)$ with respect to which $L(\lambda)$ decomposes into a direct sum of 1-dimensional $H$-modules. Let $\rho$ be the representation of $L$ afforded by this basis. Consider the function $P: L \rightarrow \mathbb{C}$ given by

$$
P(x)=\operatorname{tr}\left((\rho(x))^{m}\right) \quad x \in L
$$

We claim that $P \in P^{m}(L)$. For let $b_{1}, \ldots, b_{n}$ be a basis of $L$ and let

$$
x=\xi_{1} b_{1}+\cdots+\xi_{n} b_{n} \quad \xi_{i} \in \mathbb{C}
$$

Then we have

$$
\begin{aligned}
& \rho(x)=\sum_{i} \xi_{i} \rho\left(b_{i}\right) \\
& (\rho(x))^{m}=\sum_{i_{1}, \ldots, i_{m}} \xi_{i_{1}} \ldots \xi_{i_{m}} \rho\left(b_{i_{1}}\right) \ldots \rho\left(b_{i_{m}}\right) \\
& \operatorname{trace}(\rho(x))^{m}=\sum_{i_{1}, \ldots, i_{m}} \operatorname{tr}\left(\rho\left(b_{i_{1}}\right) \ldots \rho\left(b_{i_{m}}\right)\right) \xi_{i_{1}} \ldots \xi_{i_{m}}
\end{aligned}
$$

This is evidently a polynomial function on $L$ which is homogeneous of degree $m$. Thus $P \in P^{m}(L)$.

We wish to show that $P$ is an invariant polynomial function, that is $P \in\left(P^{m}(L)\right)^{G}$. We shall make use of the isomorphism

$$
P^{m}(L) \rightarrow T^{m}(L)_{\mathrm{sym}}^{*}
$$

obtained in Lemma 11.7. The element $f \in T^{m}(L)_{\text {sym }}^{*}$ corresponding to $P \in P^{m}(L)$ is given by

$$
f\left(x_{1} \otimes \cdots \otimes x_{m}\right)=\frac{1}{m!} \sum_{\sigma \in S_{m}} \operatorname{tr}\left(\rho\left(x_{\sigma(1)}\right) \ldots \rho\left(x_{\sigma(m)}\right)\right) .
$$

For $f$ certainly lies in $T^{m}(L)_{\text {sym }}^{*}$ and

$$
f(x \otimes \cdots \otimes x)=\operatorname{tr}\left(\rho(x)^{m}\right)
$$

Lemma 11.6 now shows that $f$ corresponds to $P$.
We recall that $T^{m}(L)$ may be regarded as an $L$-module under the action

$$
x \cdot\left(x_{1} \otimes \cdots \otimes x_{m}\right)=\sum_{i} x_{1} \otimes \cdots \otimes\left[x x_{i}\right] \otimes \cdots \otimes x_{m}
$$

Its dual space $T^{m}(L)^{*}$ then becomes an $L$-module under the action

$$
(x f)\left(x_{1} \otimes \cdots \otimes x_{m}\right)=-f\left(x\left(x_{1} \otimes \cdots \otimes x_{m}\right)\right)
$$

for $x \in L, f \in T^{m}(L)^{*}$.
We now consider $x f$ where $f \in T^{m}(L)_{\text {sym }}^{*}$ is the function defined above. We have

$$
\begin{aligned}
(x f)\left(x_{1} \otimes \cdots \otimes x_{m}\right)= & -\sum_{i} f\left(x_{1} \otimes \cdots \otimes\left[x x_{i}\right] \otimes \cdots \otimes x_{m}\right) \\
= & -\frac{1}{m!} \sum_{\sigma \in S_{m}} \sum_{i} \operatorname{tr}\left(\rho\left(x_{\sigma(1)}\right) \ldots \rho\left(\left[x x_{\sigma(i)}\right]\right) \ldots \rho\left(x_{\sigma(m)}\right)\right) \\
= & -\frac{1}{m!} \sum_{\sigma \in S_{m}} \sum_{i} \operatorname{tr}\left(\rho\left(x_{\sigma(1)}\right) \ldots \rho(x) \rho\left(x_{\sigma(i)}\right) \ldots \rho\left(x_{\sigma(m)}\right)\right) \\
& +\frac{1}{m!} \sum_{\sigma \in S_{m}} \sum_{i} \operatorname{tr}\left(\rho\left(x_{\sigma(1)}\right) \ldots \rho\left(x_{\sigma(i)}\right) \rho(x) \ldots \rho\left(x_{\sigma(m)}\right)\right) .
\end{aligned}
$$

All the terms in these expressions cancel except those for which $\rho(x)$ occurs at the beginning or the end of the product. Thus we have

$$
\begin{aligned}
&(x f)\left(x_{1} \otimes \cdots \otimes x_{m}\right)= \frac{1}{m!} \sum_{\sigma \in S_{m}}\left(\operatorname{tr}\left(\rho\left(x_{\sigma(1)}\right) \ldots \rho\left(x_{\sigma(m)}\right) \rho(x)\right)\right. \\
&\left.\quad-\operatorname{tr}\left(\rho(x) \rho\left(x_{\sigma(1)}\right) \ldots \rho\left(x_{\sigma(m))}\right)\right)\right) \\
&=\quad 0 \quad \text { since } \quad \operatorname{tr}(A B)=\operatorname{tr}(B A) .
\end{aligned}
$$

Thus $x f=0$ for all $x \in L$.

We now compare the $L$-action on $T^{m}(L)_{\text {sym }}^{*}$ with the $G$-action.
Let $x$ be an element of $L$ such that $\operatorname{ad} x$ is nilpotent. Then $\exp \operatorname{ad} x \in G$ and $G$ is generated by all such elements. Let

$$
\tau(x): T^{m}(L)_{\mathrm{sym}}^{*} \rightarrow T^{m}(L)_{\mathrm{sym}}^{*}
$$

be the linear map given by

$$
\tau(x) f^{\prime}=x f^{\prime}
$$

Then we have

$$
\tau(x) f\left(x_{1} \otimes \cdots \otimes x_{m}\right)=\sum_{i} f\left(x_{1} \otimes \cdots \otimes \operatorname{ad}(-x) \cdot x_{i} \otimes \cdots \otimes x_{m}\right)
$$

Thus

$$
\left(\frac{\tau(x)^{k}}{k!} f\right)\left(x_{1} \otimes \cdots \otimes x_{m}\right)=\sum_{\substack{i_{1}, \ldots, i_{m} \\ i_{1}+\cdots+i_{m}=k}} f\left(\frac{(\mathrm{ad}-x)^{i_{1}}}{i_{1}!} x_{1} \otimes \cdots \otimes \frac{(\mathrm{ad}-x)^{i_{m}}}{i_{m}!} x_{m}\right)
$$

Since ad $x$ is nilpotent the right-hand side is 0 for $k$ sufficiently large. Hence

$$
\begin{aligned}
(\exp \tau(x) \cdot f)\left(x_{1} \otimes \cdots \otimes x_{m}\right) & =\sum_{i_{1}, \ldots, i_{m}} f\left(\frac{(\operatorname{ad}-x)^{i_{1}}}{i_{1}!} x_{1} \otimes \cdots \otimes \frac{(\operatorname{ad}-x)^{i_{m}}}{i_{m}!} x_{m}\right) \\
& =f\left(\exp \operatorname{ad}-x \cdot x_{1} \otimes \cdots \otimes \operatorname{exp~ad}-x \cdot x_{m}\right) \\
& =(\operatorname{exp~ad}-x \cdot f)\left(x_{1} \otimes \cdots \otimes x_{m}\right)
\end{aligned}
$$

Thus we see that

$$
\exp \text { ad }-x \cdot f=\exp \tau(x) \cdot f
$$

Now we have shown that $x f=0$, hence $\tau(x) f=0$. Thus $\exp \tau(x) \cdot f=f$. It follows that

$$
\exp \text { ad }-x \cdot f=f
$$

Since this holds for all $x \in L$ with ad $x$ nilpotent we deduce that

$$
f \in\left(T^{m}(L)_{\mathrm{sym}}^{*}\right)^{G}
$$

By Lemma 11.7 it follows that $P \in P^{m}(L)^{G}$.

The restriction $\psi(P)$ therefore lies in $P^{m}(H)^{W}$. Let $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{k}$ be the weights of $L(\lambda)$ with $\lambda_{1}=\lambda$. Then we have

$$
\begin{aligned}
& \rho(x)=\left(\begin{array}{ccc}
\lambda_{1}(x) & & \\
& \cdot & \\
& \cdot & \\
& \cdot & \\
O & & \lambda_{k}(x)
\end{array}\right) \quad x \in H \\
& \rho(x)^{m}=\left(\begin{array}{ccc}
\lambda_{1}(x)^{m} & & O \\
& \cdot & \\
& \cdot & \\
O & & \lambda_{k}(x)^{m}
\end{array}\right) \\
& \operatorname{tr} \rho(x)^{m}=\lambda_{1}^{m}(x)+\cdots+\lambda_{k}^{m}(x) .
\end{aligned}
$$

Hence $\psi(P)=\lambda_{1}^{m}+\cdots+\lambda_{k}^{m}$.
We shall show that polynomial functions of this kind span $P^{m}(H)^{W}$. In the first place we know that $H^{*}$ is spanned by the lattice $X$ of integral weights. It follows that $P^{m}(H)$ is spanned by the set of monomials of degree $m$ in the integral weights. However, it is well known that the process of polarisation can be used to express such a monomial as a linear combination of $m$ th powers. (For example the formula

$$
\lambda_{1} \lambda_{2}=\frac{1}{2}\left(\lambda_{1}+\lambda_{2}\right)^{2}-\frac{1}{2} \lambda_{1}^{2}-\frac{1}{2} \lambda_{2}^{2}
$$

expresses the monomial $\lambda_{1} \lambda_{2}$ as a linear combination of squares.) Thus the elements $\lambda^{m}$ for $\lambda \in X$ span $P^{m}(H)$. It follows that every $W$-invariant element of $P^{m}(H)$ is a linear combination of elements of form

$$
\sum_{w \in W} w(\lambda)^{m} \quad \lambda \in X
$$

Since each $W$-orbit of integral weights contains a dominant integral weight we see that elements of form

$$
\sum_{w \in W} w(\lambda)^{m} \quad \lambda \in X^{+}
$$

$\operatorname{span} P^{m}(H)$.
Now we have $\psi(P)=\lambda_{1}^{m}+\cdots+\lambda_{k}^{m}$ where $\lambda_{1}=\lambda . \lambda$ appears with multiplicity 1 in the set $\left\{\lambda_{1}, \ldots, \lambda_{k}\right\}$ and each $w(\lambda)$ also appears in this set. Moreover this set is $W$-invariant, so is a union of $W$-orbits.

It follows from these facts that $\psi(P)=\sum_{w \in W} w(\lambda)^{m}+$ a linear combination of terms $\sum_{w \in W} w(\mu)^{m}$ for $\mu \in X^{+}$with $\mu \prec \lambda$. There are only finitely many
weights $\mu \in X^{+}$with $\mu \prec \lambda$. Therefore we may invert these equations and express $\sum_{w \in W} w(\lambda)^{m}$ as a linear combination of functions of the form $\psi(P)$ coming from representations with highest weight $\mu \prec \lambda$. Thus $P^{m}(H)^{W}$ is spanned by functions of the form $\psi(P)$. Hence $P^{m}(H)^{W}$ lies in the image of $\psi$. Since this is true for all $m$ the image of $\psi$ must be the whole of $P(H)^{W}$. Thus $\psi$ is surjective.

We therefore have an isomorphism of algebras

$$
\psi: P(L)^{G} \rightarrow P(H)^{W}
$$

### 11.3 The structure of the ring of polynomial invariants

In this section we shall prove a theorem of Chevalley which shows that the ring $P(H)^{W}$ of $W$-invariant polynomials on $H$ is isomorphic to a polynomial ring in $l$ variables over $\mathbb{C}$.

We write $I=P(H)^{W}$ and define $\theta: P(H) \rightarrow P(H)$ to be the operation of averaging over $W$. Thus

$$
\theta(P)=\frac{1}{|W|} \sum_{w \in W} w(P)
$$

It is clear that $\theta(P(H))=I$, that $\theta$ acts as the identity on $I$, and that $\theta^{2}=\theta$, i.e. $\theta$ is idempotent.

Let $P(H)^{+}$be the set of polynomial functions with constant term 0 , and let $I^{+}=I \cap P(H)^{+}$. Let $P(H) I^{+}$be the ideal of $P(H)$ generated by $I^{+}$. The elements of $P(H) I^{+}$have form

$$
P_{1} J_{1}+\cdots+P_{k} J_{k}
$$

with $P_{i} \in P(H), J_{i} \in I^{+}$.

Lemma 11.11 Suppose $J_{1}, \ldots, J_{k}$ are elements of $I$ such that $J_{1}$ does not lie in the ideal of $I$ generated by $J_{2}, \ldots, J_{k}$. Let $P_{1}, P_{2}, \ldots, P_{k} \in P(H)$ be homogeneous polynomials such that

$$
P_{1} J_{1}+P_{2} J_{2}+\cdots+P_{k} J_{k}=O
$$

Then $P_{1} \in P(H) I^{+}$.

Proof. We shall show that $J_{1}$ does not lie in the ideal of $P(H)$ generated by $J_{2}, \ldots, J_{k}$. Suppose this were false. Then we have

$$
J_{1}=Q_{2} J_{2}+\cdots+Q_{k} J_{k} \quad \text { with } Q_{i} \in P(H)
$$

Applying $w \in W$ we obtain

$$
J_{1}=w\left(Q_{2}\right) J_{2}+\cdots+w\left(Q_{k}\right) J_{k}
$$

and therefore

$$
J_{1}=\theta\left(Q_{2}\right) J_{2}+\cdots+\theta\left(Q_{k}\right) J_{k}
$$

However, $\theta\left(Q_{i}\right) \in I$ and so $J_{1}$ lies in the ideal of $I$ generated by $J_{2}, \ldots, J_{k}$. This gives the required contradiction.

We now show that $P_{1} \in P(H) I^{+}$by induction on the degree of the homogeneous polynomial $P_{1}$.

If $\operatorname{deg} P_{1}=0$ then $P_{1}$ is constant. Since $P_{1} J_{1}+\cdots+P_{k} J_{k}=O$ and $J_{1}$ is not in the ideal of $P(H)$ generated by $J_{2}, \ldots, J_{k}$ this implies that $P_{1}=O$. Thus $P_{1} \in P(H) I^{+}$in this case.

Now suppose $\operatorname{deg} P_{1}>0$. We recall that $W$ is generated by its fundamental reflections $s_{1}, \ldots, s_{l}$. In the $W$-action on $H$ each $s_{j}$ has a fixed point set which is a hyperplane in $H$ given by an equation $H_{j}=O$ where $H_{j} \in P(H)$ is a homogeneous polynomial of degree 1 . We have

$$
\left(s_{j}\left(P_{i}\right)\right) x=P_{i}\left(s_{j}(x)\right)=P_{i}(x)
$$

where $H_{j}(x)=0$. Thus the polynomial $s_{j}\left(P_{i}\right)-P_{i}$ vanishes at all $x \in H$ for which $H_{j}$ vanishes. It follows that

$$
s_{j}\left(P_{i}\right)-P_{i}=H_{j} \bar{P}_{i}
$$

for some $\bar{P}_{i} \in P(H)$.
Since $P_{i}$ is homogeneous, $s_{j}\left(P_{i}\right)$ is also homogeneous of the same degree, hence $s_{j}\left(P_{i}\right)-P_{i}$ is homogeneous. Thus $\bar{P}_{i}$ is also homogeneous with $\operatorname{deg} \bar{P}_{i}<$ $\operatorname{deg} P_{i}$.

Now the relation

$$
P_{1} J_{1}+\cdots+P_{k} J_{k}=O
$$

implies

$$
s_{j}\left(P_{1}\right) J_{1}+\cdots+s_{j}\left(P_{k}\right) J_{k}=O
$$

and so

$$
H_{j}\left(\bar{P}_{1} J_{1}+\cdots+\bar{P}_{k} J_{k}\right)=O
$$

Since $H_{j}$ is not the zero polynomial this implies that

$$
\bar{P}_{1} J_{1}+\cdots+\bar{P}_{k} J_{k}=O
$$

Since $\operatorname{deg} \bar{P}_{1}<\operatorname{deg} P_{1}$ we may deduce by induction that $\bar{P}_{1} \in P(H) I^{+}$. Hence $s_{j}\left(P_{1}\right)-P_{1} \in P(H) I^{+}$also.

Now $P(H) I^{+}$is a $W$-submodule of $P(H)$, thus $P(H) / P(H) I^{+}$is also a $W$-module. We have

$$
s_{j}\left(P_{1}\right) \equiv P_{1} \quad \bmod P(H) I^{+}
$$

and since $W$ is generated by $s_{1}, \ldots, s_{l}$ it follows that

$$
w\left(P_{1}\right) \equiv P_{1} \quad \bmod P(H) I^{+}
$$

for all $w \in W$. Hence

$$
\theta\left(P_{1}\right) \equiv P_{1} \quad \bmod P(H) I^{+}
$$

Now $P_{1}$ is a homogeneous polynomial of positive degree, therefore $\theta\left(P_{1}\right) \in$ $I^{+}$. In particular $\theta\left(P_{1}\right) \in P(H) I^{+}$and so $P_{1} \in P(H) I^{+}$as required.

Now the ideal $P(H) I^{+}$of $P(H)$ is generated by the homogeneous elements of $I$ of positive degree. By Hilbert's basis theorem there is a finite subset of this generating set which generates $P(H) I^{+}$. Let $I_{1}, \ldots, I_{n}$ be a set of homogeneous polynomials in $I$ such that $I_{1}, \ldots, I_{n}$ generates $P(H) I^{+}$but no proper subset generates $P(H) I^{+}$.

Proposition 11.12 The polynomials $I_{1}, \ldots, I_{n}$ are algebraically independent.

Proof. Suppose the result is false. Then there is a non-zero polynomial $P$ in $n$ variables such that

$$
P\left(I_{1}, \ldots, I_{n}\right)=O
$$

We may assume, by comparing terms of a given degree, that all monomials in $I_{1}, \ldots, I_{n}$ which occur in $P$ have the same degree $d$ in $x_{1}, \ldots, x_{l}$. Let $P_{i}=\partial P / \partial I_{i}$. Then

$$
P_{i}\left(I_{1}, \ldots, I_{n}\right) \quad i=1, \ldots, n
$$

are elements of $I$ and not all the $P_{i}$ are zero.
Let $J$ be the ideal of $I$ generated by $P_{1}, P_{2}, \ldots, P_{n}$. We may choose the notation so that $P_{1}, \ldots, P_{m}$ but no proper subset generate $J$ as an ideal in $I$. Thus there exist polynomials $Q_{i, j} \in I$ such that

$$
P_{i}=\sum_{j=1}^{m} Q_{i, j} P_{j} \quad i=m+1, \ldots, n
$$

Now each $P_{i}$ is homogeneous in $x_{1}, \ldots, x_{l}$ of degree $d-\operatorname{deg} I_{i}$. Thus, by comparing terms of the same degree in $x_{1}, \ldots, x_{l}$ on both sides, we may assume that each $Q_{i, j}$ is homogeneous of degree $\operatorname{deg} P_{i}-\operatorname{deg} P_{j}$.

Now $P\left(I_{1}, \ldots, I_{n}\right)=O$ thus $\partial P / \partial x_{k}=0$ for $k=1, \ldots, l$. Hence

$$
\sum_{i=1}^{n} \frac{\partial P}{\partial I_{i}} \frac{\partial I_{i}}{\partial x_{k}}=0
$$

that is

$$
\sum_{i=1}^{n} P_{i} \partial I_{i} / \partial x_{k}=0
$$

It follows that

$$
\sum_{i=1}^{m} P_{i} \partial I_{i} / \partial x_{k}+\sum_{i=m+1}^{n} \sum_{j=1}^{m} Q_{i, j} P_{j} \partial I_{i} / \partial x_{k}=0
$$

that is

$$
\sum_{i=1}^{m} P_{i}\left(\partial I_{i} / \partial x_{k}+\sum_{j=m+1}^{n} Q_{j, i} \partial I_{j} / \partial x_{k}\right)=0
$$

We now apply Lemma 11.11. $P_{1}, \ldots, P_{m}$ are in $I$ and $P_{1}$ is not in the ideal of $I$ generated by $P_{2}, \ldots, P_{m}$. Each of the polynomials

$$
\partial I_{i} / \partial x_{k}+\sum_{j=m+1}^{n} Q_{j, i} \partial I_{j} / \partial x_{k} \quad i=1, \ldots, m
$$

is homogeneous in $x_{1}, \ldots, x_{l}$ of degree $\operatorname{deg} I_{i}-1$. For

$$
\operatorname{deg} Q_{j, i}=\operatorname{deg} P_{j}-\operatorname{deg} P_{i}=\operatorname{deg} I_{i}-\operatorname{deg} I_{j}
$$

It follows from Lemma 11.11 that

$$
\partial I_{1} / \partial x_{k}+\sum_{j=m+1}^{n} Q_{j, 1} \partial I_{j} / \partial x_{k} \in P(H) I^{+}
$$

We now multiply this polynomial by $x_{k}$ and sum over $k=1, \ldots, l$. For a homogeneous polynomial $I_{j}$ in $x_{1}, \ldots, x_{l}$ we have, by Euler's formula,

$$
\sum_{k=1}^{l} x_{k} \frac{\partial I_{j}}{\partial x_{k}}=\operatorname{deg} I_{j} \cdot I_{j}
$$

Thus we have

$$
\operatorname{deg} I_{1} \cdot I_{1}+\sum_{j=m+1}^{n} \operatorname{deg} I_{j} \cdot Q_{j, 1} I_{j}=\sum_{i=1}^{n} I_{i} R_{i}
$$

where each $R_{i} \in P(H)^{+}$. We note that all the terms on the left-hand side are homogeneous polynomials of degree deg $I_{1}$. Comparing terms of this degree on the two sides we obtain

$$
\operatorname{deg} I_{1} \cdot I_{1}+\sum_{j=m+1}^{n} \operatorname{deg} I_{j} \cdot Q_{j, 1} I_{j}=\sum_{i} I_{i} R_{i}
$$

where the sum on the right extends over a subset of $1, \ldots, n$ not including $i=1$, since $I_{1} R_{1}$ has degree greater than $\operatorname{deg} I_{1}$. It follows that $I_{1}$ is in the ideal of $P(H)$ generated by $I_{2}, \ldots, I_{n}$. However, this contradicts the definition of $I_{1}, \ldots, I_{n}$. Thus the proposition is proved.

Proposition 11.13 Every element of $I$ is a polynomial in $I_{1}, \ldots, I_{n}$.

Proof. It is sufficient to prove this for homogeneous polynomials in I. Let $J \in I$ be homogeneous. We use induction on $\operatorname{deg} J$, the result being clear if $\operatorname{deg} J=0$. Suppose $\operatorname{deg} J>0$. Then $J \in I^{+}$and in particular $J \in P(H) I^{+}$. Thus we have

$$
J=P_{1} I_{1}+\cdots+P_{n} I_{n}
$$

for certain polynomials $P_{1}, \ldots, P_{n} \in P(H)$. Since $J, I_{1}, \ldots, I_{n}$ are all homogeneous we may clearly assume that each $P_{i}$ is homogeneous also, with

$$
\operatorname{deg} P_{i}=\operatorname{deg} J-\operatorname{deg} I_{i}
$$

Then we have

$$
J=\theta\left(P_{1}\right) I_{1}+\cdots+\theta\left(P_{n}\right) I_{n}
$$

$\theta\left(P_{1}\right), \ldots, \theta\left(P_{n}\right)$ are homogeneous polynomials in $I$ of degree less than $\operatorname{deg} J$. Thus they are polynomials in $I_{1}, \ldots, I_{n}$ by induction, and so $J$ is also.

Corollary 11.14 The algebra $P(H)^{W}=\mathbb{C}\left[I_{1}, \ldots, I_{n}\right]$ is isomorphic to the polynomial ring in $n$ generators over $\mathbb{C}$.

Proof. This follows from Propositions 11.12 and 11.13.
The set $I_{1}, \ldots, I_{n}$ is called a set of basic polynomial invariants of $W$. We now determine the number of invariants in a basic set.

Proposition 11.15 The number $n$ of invariants in a basic set is equal to the dimension lof $H$.

Proof. Let $K=\mathbb{C}\left(x_{1}, \ldots, x_{l}\right)$ be the field of rational functions in $x_{1}, \ldots, x_{l}$ over $\mathbb{C}$. Also let $k=\mathbb{C}\left(I_{1}, \ldots, I_{n}\right)$ be the field of rational functions in $I_{1}, \ldots, I_{n}$ over $\mathbb{C}$. Then we have inclusions

$$
\mathbb{C} \subset k \subset K
$$

Since $x_{1}, \ldots, x_{l}$ are algebraically independent over $\mathbb{C}$ the transcendence degree of $K$ over $\mathbb{C}$ is given by

$$
\operatorname{tr} \operatorname{deg} K / \mathbb{C}=l
$$

Since $I_{1}, \ldots, I_{n}$ are algebraically independent over $\mathbb{C}$, by Proposition 11.12, the transcendence degree of $k$ over $\mathbb{C}$ is given by

$$
\operatorname{tr} \operatorname{deg} k / \mathbb{C}=n
$$

Since we have

$$
\operatorname{tr} \operatorname{deg} K / \mathbb{C}=\operatorname{tr} \operatorname{deg} k / \mathbb{C}+\operatorname{tr} \operatorname{deg} K / k
$$

we shall consider $\operatorname{tr} \operatorname{deg} K / k$. Now $K$ is generated over $k$ by $x_{1}, \ldots, x_{l}$. However, each $x_{i}$ is an algebraic element over $k$. For the polynomial

$$
\prod_{w \in W}\left(t-w\left(x_{i}\right)\right)
$$

has $x_{i}$ as a root, and its coefficients are the elementary symmetric functions in the $w\left(x_{i}\right)$ as $w$ runs over $W$. These coefficients are $W$-invariants and therefore lie in $I$. In particular this polynomial lies in $k[t]$ and so $x_{i}$ is algebraic over $k$. Thus $K$ is generated by a finite number of algebraic elements over $k$ and so

$$
\operatorname{tr} \operatorname{deg} K / k=0
$$

It follows that

$$
\operatorname{tr} \operatorname{deg} K / \mathbb{C}=\operatorname{tr} \operatorname{deg} k / \mathbb{C}
$$

that is $n=l$.
Now the set $I_{1}, \ldots, I_{l}$ of basic polynomial invariants of $W$ is not uniquely determined. We show, however, that the degrees of these polynomials are uniquely determined.

Proposition 11.16 Let $I_{1}, \ldots, I_{l}$ and $I_{1}^{\prime}, \ldots, I_{l}^{\prime}$ be two sets of basic polynomial invariants of $W$ in $P(H)$. Then we may arrange the numbering so that

$$
\operatorname{deg} I_{i}=\operatorname{deg} I_{i}^{\prime} \quad \text { for } i=1, \ldots, l
$$

Proof. Each of $I_{1}^{\prime}, \ldots, I_{l}^{\prime}$ is expressible as a polynomial in $I_{1}, \ldots, I_{l}$ and conversely. Consider the matrices

$$
\left(\partial I_{i} / \partial I_{j}^{\prime}\right) \quad\left(\partial I_{i}^{\prime} / \partial I_{j}\right)
$$

These are inverse matrices, thus the determinant

$$
\operatorname{det}\left(\partial I_{i} / \partial I_{j}^{\prime}\right)
$$

is non-zero. It follows that for some permutation $\sigma$ of $1, \ldots, l$

$$
\prod_{i=1}^{l} \frac{\partial I_{i}}{\partial I_{\sigma(i)}^{\prime}} \neq 0
$$

By renumbering $I_{1}^{\prime}, \ldots, I_{l}^{\prime}$ if necessary we may assume $\sigma$ is the identity. Thus

$$
\prod_{i=1}^{l} \frac{\partial I_{i}}{\partial I_{i}^{\prime}} \neq 0
$$

and so $\partial I_{i} / \partial I_{i}^{\prime} \neq 0$ for each $i$. This means that $I_{i}$, as a polynomial in $I_{1}^{\prime}, \ldots, I_{l}^{\prime}$, involves $I_{i}^{\prime}$ and so

$$
\operatorname{deg} I_{i} \geq \operatorname{deg} I_{i}^{\prime}
$$

This implies that

$$
\sum_{i=1}^{l} \operatorname{deg} I_{i} \geq \sum_{i=1}^{l} \operatorname{deg} I_{i}^{\prime}
$$

By symmetry we must have equality. This implies

$$
\operatorname{deg} I_{i}=\operatorname{deg} I_{i}^{\prime} \quad \text { for each } i
$$

We summarise the results of this section in the following theorem, due to C . Chevalley.

Theorem 11.17 (a) The algebra $P(H)^{W}$ of $W$-invariant polynomials on $H$ is isomorphic to a polynomial ring in $l$ variables over $\mathbb{C}$.
(b) $P(H)^{W}$ may be generated as a polynomial ring by $l$ homogeneous invariant polynomials $I_{1}, \ldots, I_{l}$.
(c) The degrees $d_{1}, \ldots, d_{l}$ of $I_{1}, \ldots, I_{l}$ are independent of the system of generators chosen.

### 11.4 The Killing isomorphisms

In the preceding sections we have investigated the algebras $P(L)^{G}$ and $P(H)^{W}$ of invariant polynomial functions on $L$ and $H$ respectively. Assuming again that the Lie algebra $L$ is semisimple we show now how to relate these algebras
to algebras $S(L)^{G}$ and $S(H)^{W}$ of invariants on the symmetric algebras of $L$ and $H$.

The action of $G$ on the Lie algebra $L$ may be extended to a $G$-action on $T(L)$ satisfying

$$
\gamma\left(x_{1} \otimes \cdots \otimes x_{m}\right)=\gamma x_{1} \otimes \cdots \otimes \gamma x_{m} \quad \gamma \in G .
$$

We then obtain an induced action on $S(L)=T(L) / I$ since $I$ is a $G$-submodule. $S(L)^{G}$ is the subalgebra of all $G$-invariant elements of $S(L)$. We shall relate this to $P(L)^{G}$ by means of the Killing form.

We recall from Theorem 4.10 that the Killing form on the semisimple Lie algebra $L$ is non-degenerate. This implies that the linear map $L \rightarrow L^{*}$ given by $x \rightarrow x^{*}$ where $x^{*}(y)=\langle x, y\rangle$ is bijective. We wish to show that this is an isomorphism of $G$-modules.

Proposition 11.18 Let $\gamma \in G$ and $x, y \in L$. Then $\langle\gamma x, \gamma y\rangle=\langle x, y\rangle$. Thus the adjoint group preserves the Killing form.

Proof. Since $G$ is generated by elements $\exp$ ad $z$ where $z \in L$ is such that ad $z$ is nilpotent, it is sufficient to show that

$$
\langle\exp \operatorname{ad} z \cdot x, \exp \text { ad } z \cdot y\rangle=\langle x, y\rangle
$$

We recall from Proposition 4.5 that

$$
\langle[x z], y\rangle=\langle x,[z y]\rangle .
$$

Thus $\langle\operatorname{ad} z \cdot x, y\rangle=\langle x, \operatorname{ad}-z \cdot y\rangle$. Iterating we obtain

$$
\left\langle(\operatorname{ad} z)^{i} x, y\right\rangle=\left\langle x,(\operatorname{ad}-z)^{i} y\right\rangle .
$$

Now we have

$$
\exp \operatorname{ad} z=1+\operatorname{ad} z+\frac{(\operatorname{ad} z)^{2}}{2!}+\cdots+\frac{(\operatorname{ad} z)^{k}}{k!}
$$

for some $k$, since $\operatorname{ad} z$ is nilpotent. Hence

$$
\langle\exp \operatorname{ad} z \cdot x, y\rangle=\langle x, \exp \operatorname{ad}-z \cdot y\rangle
$$

and so

$$
\langle\exp \operatorname{ad} z \cdot x, \exp \operatorname{ad} z \cdot y\rangle=\langle x, y\rangle
$$

Corollary 11.19 The Killing map $L \rightarrow L^{*}$ is an isomorphism of $G$-modules.

Proof. We must show that $(\gamma x)^{*}=\gamma x^{*}$ for all $\gamma \in G, x \in L$. We have

$$
(\gamma x)^{*}(y)=\langle\gamma x, y\rangle=\left\langle x, \gamma^{-1} y\right\rangle=x^{*}\left(\gamma^{-1} y\right)=\left(\gamma x^{*}\right)(y)
$$

Thus $(\gamma x)^{*}=\gamma x^{*}$ as required.
The Killing map $L \rightarrow L^{*}$ induces an isomorphism $T(L) \rightarrow T\left(L^{*}\right)$ and then an isomorphism $S(L) \rightarrow S\left(L^{*}\right)$ in an obvious way. This is again an isomorphism of $G$-modules. There is therefore an isomorphism between $S(L)^{G}$ and $S\left(L^{*}\right)^{G}$. We recall that $S\left(L^{*}\right)=P(L)$ and so obtain a Killing isomorphism of algebras $S(L) \rightarrow P(L)$ which induces a Killing isomorphism $S(L)^{G} \rightarrow P(L)^{G}$ between the subalgebras of invariants.

We now consider the action of the Weyl group $W$ on the Cartan subalgebra $H$ of $L$. We recall from Proposition 4.14 that the Killing form of $L$ remains non-degenerate on restriction to $H$. Thus the map $H \rightarrow H^{*}$ given by $x \rightarrow x^{*}$ where $x^{*}(y)=\langle x, y\rangle$ for all $y \in H$ is bijective.

Proposition 11.20 The Killing map $H \rightarrow H^{*}$ is an isomorphism of $W$-modules.

Proof. We have

$$
(w h)^{*} x=\langle w h, x\rangle=\left\langle h, w^{-1} x\right\rangle=h^{*}\left(w^{-1} x\right)=\left(w h^{*}\right) x \quad \text { for all } x \in H
$$

Hence $(w h)^{*}=w h^{*}$ as required.
The Killing isomorphism $H \rightarrow H^{*}$ induces an isomorphism $T(H) \rightarrow T\left(H^{*}\right)$ and then an isomorphism $S(H) \rightarrow S\left(H^{*}\right)$. This is again an isomorphism of $W$ modules. Since $S\left(H^{*}\right)=P(H)$ we obtain a Killing isomorphism of algebras $S(H) \rightarrow P(H)$ which induces an isomorphism $S(H)^{W} \rightarrow P(H)^{W}$ between the subalgebras of invariants.

We now consider the relation between $S(L)$ and $S(H)$. We recall that $L$ may be identified with a subspace of $S(L)$ and that $L$ has a triangular decomposition

$$
L=N^{-} \oplus H \oplus N
$$

Let $K$ be the ideal of $S(L)$ generated by $N$ and $N^{-}$. Then we have $S(L) / K$ isomorphic to $S(H)$. Let $\eta: S(L) \rightarrow S(H)$ be the natural homomorphism given in this way.

Proposition 11.21 We have a commutative diagram of algebra homomorphisms

$$
\begin{array}{cc}
S(L) \xrightarrow{\alpha} P(L) \\
\downarrow & \downarrow \\
\eta & \downarrow \\
S(H) & \underset{\beta}{\rightarrow} P(H)
\end{array}
$$

where $\alpha, \beta$ are the Killing isomorphisms, $\psi$ is restriction from $P(L)$ to $P(H)$, and $\eta$ is projection from $S(L)$ to $S(H)$.

Proof. We must show $\psi \alpha(Q)=\beta \eta(Q)$ for all $Q \in S(L)$. It is sufficient to prove this when

$$
Q=f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}
$$

where $\Phi^{+}=\left\{\beta_{1}, \ldots, \beta_{N}\right\}$.
If $r_{i}=0$ and $t_{i}=0$ for each $i$ then $\eta(Q)=Q$. Moreover $\beta(Q)=\psi \alpha(Q)$. Thus the diagram commutes.

If not all the $r_{i}$ and $t_{i}$ are 0 then $\eta(Q)=O$. Thus $\beta \eta(Q)=O$. We have

$$
\alpha(Q)=\alpha\left(f_{\beta_{1}}\right)^{r_{1}} \ldots \alpha\left(f_{\beta_{N}}\right)^{r_{N}} \alpha\left(h_{1}\right)^{s_{1}} \ldots \alpha\left(h_{l}\right)^{s_{l}} \alpha\left(e_{\beta_{1}}\right)^{t_{1}} \ldots \alpha\left(e_{\beta_{N}}\right)^{t_{N}}
$$

Therefore, for $x \in H$ we have

$$
(\alpha Q) x=\left\langle f_{\beta_{1}}, x\right\rangle^{r_{1}} \ldots\left\langle f_{\beta_{N}}, x\right\rangle^{r_{N}}\left\langle h_{1}, x\right\rangle^{s_{1}} \ldots\left\langle h_{l_{1}}, x\right\rangle^{s_{l}}\left\langle e_{\beta_{1}}, x\right\rangle^{t_{1}} \ldots\left\langle e_{\beta_{N}}, x\right\rangle^{t_{N}}
$$

This is 0 since $\left\langle N^{-}, H\right\rangle=0$ and $\langle N, H\rangle=0$, and some $r_{i}$ or $t_{i}$ is non-zero. Thus the diagram commutes in this case also.

Corollary 11.22 We have a commutative diagram of algebra isomorphisms

$$
\begin{array}{cc}
S(L)^{G} \xrightarrow{\alpha} P(L)^{G} \\
\downarrow & \downarrow \\
\eta & \downarrow \\
S(H)^{W} & \underset{\beta}{\longrightarrow} P(H)^{W}
\end{array}
$$

Proof. We have seen that the Killing isomorphisms $\alpha, \beta$ map $S(L)^{G}$ to $P(L)^{G}$ and $S(H)^{W}$ to $P(H)^{W}$, respectively. We also know from Theorem 11.10 that $\psi: P(L)^{G} \rightarrow P(H)^{W}$ is an isomorphism of algebras. Thus $\eta$ acts on $S(L)^{G}$ in the same way as $\beta^{-1} \psi \alpha$. Hence $\eta: S(L)^{G} \rightarrow S(H)^{W}$ is an algebra isomorphism.

We note by Theorem 11.17 that the four algebras $S(L)^{G}, P(L)^{G}, S(H)^{W}$, $P(H)^{W}$ are all isomorphic to the polynomial algebra $\mathbb{C}\left[z_{1}, \ldots, z_{l}\right]$.

### 11.5 The centre of the enveloping algebra

The centre $Z(L)$ of $\mathfrak{H}(L)$ is defined by

$$
Z(L)=\{z \in \mathfrak{U}(L) ; z u=u z \quad \text { for all } u \in \mathfrak{U}(L)\}
$$

Proposition 11.23 The centre $Z(L)$ acts on each Verma module $M(\lambda)$ by scalar multiplications.

Proof. Let $m_{\lambda}$ be the highest weight vector of $M(\lambda)$. Let $z \in Z(L)$ and $h \in H$. Then

$$
h\left(z m_{\lambda}\right)=z\left(h m_{\lambda}\right)=\lambda(h) z m_{\lambda} .
$$

Thus $z m_{\lambda} \in M(\lambda)_{\lambda}$. Now the $\lambda$-weight space of $M(\lambda)$ is 1 -dimensional - in fact $M(\lambda)_{\lambda}=\mathbb{C} m_{\lambda}$. Hence

$$
z m_{\lambda}=\xi m_{\lambda} \quad \text { for some } \xi \in \mathbb{C}
$$

Now let $u \in \mathfrak{H}(L)$. Then we have

$$
z\left(u m_{\lambda}\right)=u\left(z m_{\lambda}\right)=\xi u m_{\lambda} .
$$

Since $M(\lambda)=\mathfrak{l l}(L) m_{\lambda}$ we see that $z$ acts on $M(\lambda)$ as scalar multiplication by $\xi$.

We write $\chi_{\lambda}(z)=\xi$. Thus $\chi_{\lambda}: Z(L) \rightarrow \mathbb{C}$ is a 1-dimensional representation of $Z(L) . \chi_{\lambda}$ is called the central character of $M(\lambda)$. We shall show how to determine this central character.

We consider $\mathfrak{H}(L)$ as an $L$-module, as described in Section 11.1. The $L$-action on $\mathfrak{U}(L)$ is given by

$$
x \cdot u=x u-u x \quad x \in L, u \in \mathfrak{H}(L)
$$

$\mathfrak{H}(L)$ has basis

$$
f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} \quad h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} \quad e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}
$$

where $\Phi^{+}=\left\{\beta_{1}, \ldots, \beta_{N}\right\}$. If $x \in H$ we have

$$
\begin{aligned}
& x \cdot f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} \quad h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} \quad e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}=\left(-r_{1} \beta_{1}-\cdots-r_{N} \beta_{N}\right. \\
& \left.\quad+t_{1} \beta_{1}+\cdots+t_{N} \beta_{N}\right)(x) f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} \quad h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} \quad e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}} .
\end{aligned}
$$

Thus $f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} \quad h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} \quad e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}$ is a weight vector with weight $\left(t_{1}-r_{1}\right) \beta_{1}+\cdots+\left(t_{N}-r_{N}\right) \beta_{N}$.

We consider the zero weight space $\mathfrak{l}(L)_{0}$. This has basis $f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{T_{N}}$ $h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} \quad e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}$ where $\left(t_{1}-r_{1}\right) \beta_{1}+\cdots+\left(t_{N}-r_{N}\right) \beta_{N}=0$. We have

$$
\mathfrak{H}(L)_{0}=\{u \in \mathfrak{H}(L) ; x u-u x=0 \quad \text { for all } x \in H\}
$$

thus $\mathfrak{H}(L)_{0}$ is a subalgebra of $\mathfrak{H}(L)$. It is clear that $Z(L) \subset \mathfrak{H}(L)_{0}$.
Proposition 11.24 (i) $\mathfrak{l}(L) N \cap \mathfrak{U}(L)_{0}=N^{-} \mathfrak{l}(L) \cap \mathfrak{l}(L)_{0}=K$.
(ii) The subspace $K$ of (i) is a 2-sided ideal of $\mathfrak{H}(L)_{0}$.
(iii) $\mathfrak{U}(L)_{0}=K \oplus \mathfrak{U}(H)$.

Proof. (i) $\mathfrak{l}(L) N$ is spanned by the basis vectors of $\mathfrak{H}(L)$ with some $t_{i}>0$. $N^{-} \mathfrak{l}(L)$ is spanned by the basis vectors with some $r_{i}>0$. $\mathfrak{U}(L) N \cap \mathfrak{U}(L)_{0}$ is spanned by the basis vectors of $\mathfrak{H}(L)$ with $\sum t_{i} \beta_{i}=$ $\sum r_{i} \beta_{i}$ and some $t_{i}>0 . N^{-} \mathfrak{l}(L) \cap \mathfrak{l}(L)_{0}$ is spanned by the basis vectors of $\mathfrak{H}(L)$ with $\sum t_{i} \beta_{i}=\sum r_{i} \beta_{i}$ and some $r_{i}>0$. These are clearly equal.
(ii) $\mathfrak{U}(L) N \cap \mathfrak{l}(L)_{0}$ is clearly a left ideal of $\mathfrak{H}(L)_{0}$ and $N^{-} \mathfrak{l}(L) \cap \mathfrak{l}(L)_{0}$ is a right ideal of $\mathfrak{U}(L)_{0}$. Thus $K$ is a 2 -sided ideal of $\mathfrak{U t}(L)_{0}$.
(iii) $\mathfrak{U}(H)$ is spanned by the basis vectors with all $r_{i}=0$ and all $t_{i}=0$. This shows that $\mathfrak{l}(L)_{0}$ is the direct sum of its subspaces $K$ and $\mathfrak{l}(H)$.

Let $\phi: \mathfrak{U}(L)_{0} \rightarrow \mathfrak{U}(H)$ be the projection map obtained from the decomposition

$$
\mathfrak{U}(L)_{0}=K \oplus \mathfrak{l}(H) .
$$

Since $K$ is a 2-sided ideal of $\mathfrak{H}(L)_{0}, \phi$ is a homomorphism of algebras. $\phi$ is called the Harish-Chandra homomorphism.

We can now determine the central character $\chi_{\lambda}$. The weight $\lambda \in H^{*}$ determines a 1 -dimensional representation of $\mathfrak{H}(H)$, also denoted by $\lambda$.

Theorem 11.25 The central character $\chi_{\lambda}: Z(L) \rightarrow \mathbb{C}$ is given by $\chi_{\lambda}(z)=$ $\lambda(\phi(z))$ where $\phi$ is the Harish-Chandra homomorphism.

Proof. We have

$$
\mathfrak{U}(L)_{0}=\left(\mathfrak{l}(L) N \cap \mathfrak{U}(L)_{0}\right) \oplus \mathfrak{U}(H)
$$

and $Z(L) \subset \mathfrak{H}(L)_{0}$. Let $z \in Z(L)$. Then we can write

$$
z=u_{1} n_{1}+\cdots+u_{k} n_{k}+\phi(z)
$$

where $u_{i} \in \mathfrak{H}(L)$ and $n_{i} \in N$. Thus

$$
\begin{aligned}
z m_{\lambda} & =\left(u_{1} n_{1}+\cdots+u_{k} n_{k}+\phi(z)\right) m_{\lambda} \\
& =\lambda(\phi(z)) m_{\lambda}
\end{aligned}
$$

since $N m_{\lambda}=O$ and $\phi(z) m_{\lambda}=\lambda(\phi(z)) m_{\lambda}$. Thus $\chi_{\lambda}(z)=\lambda(\phi(z))$.
We have seen that the Harish-Chandra homomorphism maps $Z(L)$ into $\mathfrak{U}(H)$. Since the Lie algebra $H$ is abelian we have $\mathfrak{H}(H)=S(H)$. We shall show that by combining the Harish-Chandra homomorphism with a 'twisting homomorphism' we get a homomorphism from $Z(L)$ into $S(H)$ with very favourable properties. The twisting homomorphism $\tau: S(H) \rightarrow S(H)$ is defined as follows. We recall that $S(H)$ is a polynomial algebra over $\mathbb{C}$ with generators $h_{1}, \ldots, h_{l}$. Thus there is a unique algebra homomorphism

$$
\tau: S(H) \rightarrow S(H)
$$

such that $\tau\left(h_{i}\right)=h_{i}-1 . \tau$ is in fact an automorphism of algebras. Its inverse is given by $\tau^{-1}\left(h_{i}\right)=h_{i}+1$.

Let $\rho \in X$ be the element of the weight lattice given by

$$
\rho=\omega_{1}+\cdots+\omega_{l}
$$

Thus $\rho$ is the sum of the fundamental weights. We recall from Section 10.3 that

$$
\omega_{i}\left(h_{i}\right)=1 \quad w_{i}\left(h_{j}\right)=0 \quad \text { if } j \neq i
$$

Thus $\rho\left(h_{i}\right)=1$ for each $i=1, \ldots, l$.
Now any element $\lambda \in H^{*}$ extends to a 1-dimensional representation of $S(H)$. $\lambda-\rho$ is also a 1 -dimensional representation of $S(H)$. We have

$$
\lambda \tau\left(h_{i}\right)=\lambda\left(h_{i}-1\right)=(\lambda-\rho) h_{i} .
$$

Since $\lambda \tau$ and $\lambda-\rho$ are 1 -dimensional representations of $S(H)$ and the $h_{i}$ generate $S(H)$ we have

$$
\lambda \tau(Q)=(\lambda-\rho)(Q) \quad \text { for all } Q \in S(H)
$$

The homomorphism

$$
\tau \phi: Z(L) \rightarrow S(H)
$$

is called the twisted Harish-Chandra homomorphism. We wish to show that the image of $Z(L)$ under the twisted Harish-Chandra homomorphism lies in $S(H)^{W}$. To do so we first need a result on Verma modules.

Proposition 11.26 Let $\lambda \in H^{*}$ and $M(\lambda)$ be the corresponding Verma module with highest weight vector $m_{\lambda}$. Suppose $(\lambda+\rho)\left(h_{i}\right) \in \mathbb{Z}$ and $(\lambda+\rho)\left(h_{i}\right)>0$ for some i. Let

$$
v=f_{i}^{(\lambda+\rho)\left(h_{i}\right)} m_{\lambda}
$$

Then the submodule of $M(\lambda)$ generated by $v$ is isomorphic to $M(\mu)$ where

$$
\mu+\rho=s_{i}(\lambda+\rho) .
$$

Proof. We recall from Theorem 10.6 that there is an isomorphism of $\mathfrak{U}\left(N^{-}\right)$-modules between $\mathfrak{U}\left(N^{-}\right)$and $M(\lambda)$ given by $u \rightarrow u m_{\lambda}$. Since $f_{i}^{(\lambda+\rho)\left(h_{i}\right)} \neq 0$ in $\mathfrak{U}\left(N^{-}\right)$we see that $v \neq 0$ in $M(\lambda)$. Since $m_{\lambda} \in M(\lambda)_{\lambda}$ we have $v \in M(\lambda)_{\mu}$ where

$$
\mu=\lambda-(\lambda+\rho)\left(h_{i}\right) \alpha_{i} .
$$

Thus we have

$$
\mu+\rho=(\lambda+\rho)-(\lambda+\rho)\left(h_{i}\right) \alpha_{i}=s_{i}(\lambda+\rho) .
$$

We shall show that $N v=O$. It is sufficient to show that $e_{j} v=0$ for $j=1, \ldots, l$. If $j \neq i$ we have

$$
e_{j} v=e_{j} f_{i}^{(\lambda+\rho)\left(h_{i}\right)} m_{\lambda}=f_{i}^{(\lambda+\rho)\left(h_{i}\right)} e_{j} m_{\lambda}=0
$$

If $j=i$ we have

$$
\begin{aligned}
e_{i} v & =e_{i} f_{i}^{(\lambda+\rho)\left(h_{i}\right)} m_{\lambda} \\
& \left.=f_{i}^{(\lambda+\rho)\left(h_{i}\right)} e_{i}+(\lambda+\rho)\left(h_{i}\right) f_{i}^{(\lambda+\rho)\left(h_{i}\right)-1}\left(h_{i}-(\lambda+\rho)\left(h_{i}\right)-1\right)\right) m_{\lambda} \\
& =(\lambda+\rho)\left(h_{i}\right) f_{i}^{(\lambda+\rho)\left(h_{i}\right)-1}\left(\lambda\left(h_{i}\right)-\lambda\left(h_{i}\right)-1+1\right) m_{\lambda}=0 .
\end{aligned}
$$

Thus $N v=O$.
Let $V$ be the submodule of $M(\lambda)$ generated by $v$. Since $N v=O$ and $h_{i} v=$ $\mu\left(h_{i}\right) v$ for $i=1, \ldots, l$, there is a surjective homomorphism of $\mathfrak{H}(L)$ - modules from $M(\mu)$ into $V$ given by

$$
u m_{\mu} \rightarrow u v \quad u \in \mathfrak{U}\left(N^{-}\right) .
$$

(See Proposition 10.13.) We consider the kernel of this homomorphism. Let $u \in \mathfrak{U}\left(N^{-}\right)$be such that $u v=0$. Then

$$
u f_{i}^{(\lambda+\rho)\left(h_{i}\right)} m_{\lambda}=0
$$

Since $u f_{i}^{(\lambda+\rho)\left(h_{i}\right)} \in \mathfrak{H}\left(N^{-}\right)$this implies that $u f_{i}^{(\lambda+\rho)\left(h_{i}\right)}=0$. Since $f_{i}^{(\lambda+\rho)\left(h_{i}\right)} \neq 0$ and $\mathfrak{U}(L)$ has no zero-divisors we have $u=0$. Thus our homomorphism is an isomorphism and so $V$ is isomorphic to $M(\mu)$

Proposition 11.27 The twisted Harish-Chandra homomorphism $\tau \phi$ maps $Z(L)$ into $S(H)^{W}$.

Proof. We must show that $\tau \phi(z) \in S(H)^{W}$ for all $z \in Z(L)$. Since $W$ is generated by $s_{1}, \ldots, s_{l}$ it will be sufficient to show that

$$
s_{i}(\tau \phi(z))=\tau \phi(z)
$$

Since $S(H)=P\left(H^{*}\right)$ it will be sufficient to show these elements take the same value for all $\lambda \in H^{*}$, i.e. that

$$
\lambda\left(s_{i}(\tau \phi(z))\right)=\lambda(\tau \phi(z)) \quad \text { for all } \lambda \in H^{*} .
$$

In fact it will be sufficient to prove this for elements of $H^{*}$ of the form $\lambda+\rho$ where $\lambda \in X^{+}$is dominant and integral. For such weights form a dense subset of $H^{*}$ in the Zariski topology, for which the closed sets are the algebraic sets.

Thus suppose $\lambda \in X^{+}$. Then we have

$$
(\lambda+\rho)(\tau(\phi(z)))=\lambda(\phi(z))=\chi_{\lambda}(z)
$$

using Theorem 11.25 and the definition of $\tau$. Similarly we have

$$
(\lambda+\rho)\left(s_{i}(\tau(\phi(z)))\right)=(\mu+\rho)(\tau(\phi(z)))=\mu(\phi(z))=\chi_{\mu}(z)
$$

where $s_{i}(\lambda+\rho)=\mu+\rho$.
We now apply Proposition 11.26. Since $\lambda \in X^{+}$we have $\lambda\left(h_{i}\right) \geq 0$, so $(\lambda+\rho)\left(h_{i}\right)>0$. Thus the Verma module $M(\lambda)$ contains a submodule isomorphic to $M(\mu)$. Now $z \in Z(L)$ acts on $M(\lambda)$ as scalar multiplication by $\chi_{\lambda}(z)$ and on $M(\mu)$ as scalar multiplication by $\chi_{\mu}(z)$. Since $M(\mu)$ is isomorphic to a submodule of $M(\lambda)$ we must have

$$
\chi_{\lambda}(z)=\chi_{\mu}(z)
$$

Thus

$$
(\lambda+\rho)(\tau(\phi(z)))=(\lambda+\rho)\left(s_{i}(\tau(\phi(z)))\right)
$$

and hence

$$
\tau(\phi(z))=s_{i}(\tau(\phi(z)))
$$

Thus $\tau \phi(z) \in S(H)^{W}$ as required.
In fact we shall show that the twisted Harish-Chandra map

$$
\tau \phi: Z(L) \rightarrow S(H)^{W}
$$

is an isomorphism of algebras.

To see this we first recall the operation $\theta: S(L) \rightarrow \mathfrak{U}(L)$ of symmetrisation which was shown in Proposition 11.4 to be an isomorphism of $L$-modules. Now the adjoint group $G$ acts on both $S(L)$ and $\mathfrak{H}(L)$. For the $G$-action on $L$ can be extended to a $G$-action on $T(L)$ as described in Section 11.4 and these induce $G$-actions on the quotients $S(L)$ and $\mathfrak{H}(L)$. Suppose $x \in L$ is such that ad $x$ is nilpotent. Then $\exp$ ad $x \in G$. Let $x$ induce the linear maps $\alpha(x)$ on $S(L)$ and $\beta(x)$ on $\mathfrak{H}(L)$. The definition of the $G$-actions then shows that exp ad $x$ acts as $\exp \alpha(x)$ on $S(L)$ and as $\exp \beta(x)$ on $\mathfrak{U}(L)$. Since $\theta$ is an isomorphism of $L$-modules we have

$$
\theta \alpha(x)=\beta(x) \theta
$$

It follows that

$$
\theta \frac{\alpha(x)^{i}}{i!}=\frac{\beta(x)^{i}}{i!} \theta \quad \text { for all } i
$$

and therefore that

$$
\theta \exp \alpha(x)=\exp \beta(x) \theta
$$

(Note that both $\alpha(x)$ and $\beta(x)$ are nilpotent.) Since $G$ is generated by such elements $\exp \operatorname{ad} x$ it follows that $\theta$ is an isomorphism of $G$-modules. We deduce that $\theta$ restricts to an isomorphism between $S(L)^{G}$ and $\mathfrak{U}(L)^{G}$.

Proposition 11.28 $\mathfrak{U}(L)^{G}=Z(L)$.
Proof. We first note that $Z(L) \subset \mathfrak{H}(L)^{G}$. Let $z \in Z(L)$. Let $x \in L$ be such that $\operatorname{ad} x$ is nilpotent. Thus $\exp$ ad $x \in G$. Since $z \in Z(L)$ we have

$$
x \cdot z=x z-z x=0
$$

Hence $\beta(x) z=0$. Thus

$$
\exp \text { ad } x \cdot z=\exp \beta(x) \cdot z=\left(1+\beta(x)+\frac{\beta(x)^{2}}{2!}+\cdots\right) z=z
$$

Thus $z$ is invariant under $\exp$ ad $x$. Since such elements generate $G$ we have $z \in \mathfrak{U}(L)^{G}$.

Conversely we show that $\mathfrak{H}(L)^{G} \subset Z(L)$. Let $u \in \mathfrak{U}(L)^{G}$. Then exp ad $x \cdot u=$ $u$ for all $x \in L$ with ad $x$ nilpotent. Suppose $(\operatorname{ad} x)^{t} \neq 0$ but $(\operatorname{ad} x)^{t+1}=0$. We choose elements $\xi_{1}, \ldots, \xi_{t+1} \in \mathbb{C}$ which are all distinct. Then $\operatorname{ad}\left(\xi_{i} x\right)$ is also nilpotent and

$$
\begin{aligned}
\exp \operatorname{ad}\left(\xi_{i} x\right) & =1+\operatorname{ad}\left(\xi_{i} x\right)+\cdots+\frac{1}{t!}\left(\operatorname{ad}\left(\xi_{i} x\right)\right)^{t} \\
& =1+\xi_{i}(\operatorname{ad} x)+\cdots+\frac{\xi_{i}^{t}}{t!}(\operatorname{ad} x)^{t}
\end{aligned}
$$

Now the determinant

$$
\left|\begin{array}{ccccc}
1 & \xi_{1} & \frac{\xi_{1}^{2}}{2!} & \ldots & \frac{\xi_{1}^{t}}{t!} \\
1 & \xi_{2} & \frac{\xi_{2}^{2}}{2!} & \ldots & \frac{\xi_{2}^{t}}{t!} \\
\cdot & \cdot & \cdot & \ldots & \cdot \\
\cdot & \cdot & \cdot & \cdots & \cdot \\
\cdot & \cdot & \cdot & \cdots & \cdot \\
1 & \xi_{t+1} & \frac{\xi_{t+1}^{2}}{2!} & \cdots & \frac{\xi_{t+1}^{t}}{t!}
\end{array}\right|=\frac{1}{2!3!\ldots t!} \prod_{i<j}\left(\xi_{i}-\xi_{j}\right)
$$

is non-zero. Thus the vector $(0,1,0, \ldots, 0)$ is a linear combination of the rows of the determinant. Thus there exist $\eta_{1}, \ldots, \eta_{t+1} \in \mathbb{C}$ such that

$$
\operatorname{ad} x=\eta_{1} \exp \text { ad }\left(\xi_{1} x\right)+\cdots+\eta_{t+1} \exp \operatorname{ad}\left(\xi_{t+1} x\right)
$$

So ad $x \cdot u=\left(\eta_{1}+\cdots+\eta_{t+1}\right) u$. Since ad $x$ acts nilpotently on $u$ it follows that $\eta_{1}+\cdots+\eta_{t+1}=0$ and that ad $x \cdot u=0$. This means that $x u-u x=0$. This holds for all $x \in L$ with ad $x$ nilpotent, in particular for $x=e_{i}$ and $x=f_{i}$. However, $e_{1}, \ldots, e_{l}, f_{1}, \ldots, f_{l}$ generate $\mathfrak{H}(L)$, together with 1 . It follows that $x u-u x=0$ for all $x \in \mathfrak{U}(L)$, that is $u \in Z(L)$.

Thus the operation $\theta$ of symmetrisation gives an isomorphism of vector spaces

$$
\theta: S(L)^{G} \rightarrow Z(L)
$$

Now we also have an isomorphism of algebras

$$
\eta: S(L)^{G} \rightarrow S(H)^{W}
$$

given in Corollary 11.22. Combining these maps we obtain an isomorphism of vector spaces

$$
\eta \theta^{-1}: Z(L) \rightarrow S(H)^{W}
$$

Thus we have two maps $\eta \theta^{-1}$ and $\tau \phi$ from $Z(L)$ into $S(H)^{W}$. The first is an isomorphism of vector spaces and the second a homomorphism of algebras. We shall compare these maps, using the structure of $Z(L)$ and $S(H)$ as filtered algebras.

We recall from Section 11.1 that $\mathfrak{H}(L)$ may be regarded as a filtered algebra with filtration

$$
\mathfrak{H}_{0}(L) \subset \mathfrak{H}_{1}(L) \subset \mathfrak{H}_{2}(L) \subset \cdots
$$

We define $Z_{i}(L)=Z(L) \cap \mathfrak{U}_{i}(L)$. This makes $Z(L)$ into a filtered algebra. $S(H)$ also has a natural structure as a filtered algebra, where $S_{i}(H)$ is the subspace of $S(H)$ generated by all products $a_{1} a_{2} \ldots a_{j}, j \leq i$, where $a_{k} \in H$. We also define

$$
\left(S(H)^{W}\right)_{i}=S(H)^{W} \cap S_{i}(H) .
$$

This makes $S(H)^{W}$ into a filtered algebra.
We shall make use of the following lemma on filtered and graded algebras.

Lemma 11.29 Let $A=\bigcup_{i \geq 0} A_{i}$ and $B=\bigcup_{i \geq 0} B_{i}$ be filtered algebras with

$$
A_{0} \subset A_{1} \subset A_{2} \subset \cdots
$$

and

$$
B_{0} \subset B_{1} \subset B_{2} \subset \cdots
$$

Let

$$
\operatorname{gr} A=A_{0} \oplus A_{1} / A_{0} \oplus A_{2} / A_{1} \oplus \cdots
$$

and

$$
\operatorname{gr} B=B_{0} \oplus B_{1} / B_{0} \oplus B_{2} / B_{1} \oplus \cdots
$$

be the corresponding graded algebras. Let $\alpha: A \rightarrow B$ be a linear map such that $\alpha\left(A_{i}\right) \subset B_{i}$ for each $i$. Then:
(a) There is a linear map gr $\alpha$ : $\operatorname{gr} A \rightarrow \operatorname{gr} B$ satisfying $\operatorname{gr} \alpha\left(A_{i-1}+a_{i}\right)=B_{i-1}+$ $\alpha\left(a_{i}\right)$ for $a_{i} \in A_{i}$.
(b) If $\alpha\left(A_{i}\right)=B_{i}$ for each $i$ and $\alpha$ is bijective then $\operatorname{gr} \alpha$ is bijective.
(c) If $\operatorname{gr} \alpha$ is bijective then $\alpha$ is bijective.

Proof. (a) We must show that $\operatorname{gr} \alpha: A_{i} / A_{i-1} \rightarrow B_{i} / B_{i-1}$ is well defined. Suppose $A_{i-1}+a_{i}=A_{i-1}+a_{i}^{\prime}$ where $a_{i}, a_{i}^{\prime} \in A_{i}$. Then $a_{i}-a_{i}^{\prime} \in A_{i-1}$, so $\alpha\left(a_{i}-a_{i}^{\prime}\right) \in B_{i-1}$. Thus $B_{i-1}+\alpha\left(a_{i}\right)=B_{i-1}+\alpha\left(a_{i}^{\prime}\right)$ and so $\operatorname{gr} \alpha$ is well defined.
(b) Suppose now that $\alpha\left(A_{i}\right)=B_{i}$ for each $i$ and that $\alpha$ is bijective. Then the induced map $\operatorname{gr} \alpha: A_{i} / A_{i-1} \rightarrow B_{i} / B_{i-1}$ is bijective. It follows that $\operatorname{gr} \alpha: \operatorname{gr} A \rightarrow \operatorname{gr} B$ is bijective.
(c) Suppose conversely that $\operatorname{gr} \alpha: \operatorname{gr} A \rightarrow \operatorname{gr} B$ is bijective. This implies that

$$
\operatorname{gr} \alpha: A_{i} / A_{i-1} \rightarrow B_{i} / B_{i-1}
$$

is bijective for each $i$. We show first that $\alpha$ is surjective. $B_{0}$ lies in the image of $\alpha$ since $\alpha: A_{0} \rightarrow B_{0}$ agrees with gr $\alpha: A_{0} \rightarrow B_{0}$. Assume by
induction that $B_{i-1}$ lies in the image of $\alpha$. Let $b_{i} \in B_{i}$. Then there exists $a_{i} \in A_{i}$ such that

$$
B_{i-1}+\alpha\left(a_{i}\right)=B_{i-1}+b_{i}
$$

Thus $b_{i}-\alpha\left(a_{i}\right) \in B_{i-1}$. Hence $b_{i}-\alpha\left(a_{i}\right)$ lies in the image of $\alpha$, thus $b_{i}$ does also. Thus $\alpha$ is surjective.

Now let $a \in \operatorname{ker} \alpha$. If $a \in A_{0}$ then $a=0$ since $\alpha$ agrees with gr $\alpha$ on $A_{0}$. Otherwise there exists $i>0$ such that $a \in A_{i}$ but $a \notin A_{i-1}$. But then $A_{i-1}+a \neq 0$ whereas gr $\alpha\left(A_{i-1}+a\right)=0$, a contradiction. Hence ker $\alpha=O$ and so $\alpha$ is bijective.

Theorem 11.30 The twisted Harish-Chandra map $\tau \phi$ gives an isomorphism of algebras $Z(L) \rightarrow S(H)^{W}$.

Proof. We have maps $\tau \phi: Z(L) \rightarrow S(H)^{W}$ and $\eta \theta^{-1}: Z(L) \rightarrow S(H)^{W}$. Those induce maps

$$
\begin{aligned}
& \operatorname{gr}(\tau \phi): \operatorname{gr} Z(L) \rightarrow \operatorname{gr} S(H)^{W} \\
& \operatorname{gr}\left(\eta \theta^{-1}\right): \operatorname{gr} Z(L) \rightarrow \operatorname{gr} S(H)^{W}
\end{aligned}
$$

We shall show that $\operatorname{gr}(\tau \phi)=\operatorname{gr}\left(\eta \theta^{-1}\right)$. Let $z \in Z(L)$. Then there exists $d$ such that $z \in Z_{d}(L)$ but $z \notin Z_{d-1}(L)$. Then $z$ has the form

$$
z=\sum_{\sum r_{i}+\sum s_{i}+\sum t_{i} \leq d} \xi(\underline{\mathrm{r}}, \underline{\mathrm{~s}}, \underline{\mathrm{t}}) f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} h_{1}^{s_{l}} \ldots h_{l}^{s_{l}} e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}}
$$

where $\Phi^{+}=\left\{\beta_{1}, \ldots, \beta_{N}\right\}$ and $\xi(\underline{\mathrm{r}}, \underline{\mathrm{s}}, \underline{\mathrm{t}}) \in \mathbb{C}$. Then

$$
\begin{aligned}
& \phi(z)=\sum_{\sum s_{i} \leq d} \xi(\underline{\mathbf{o}}, \underline{\mathrm{~s}}, \underline{\mathrm{o}}) h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} \\
& \tau \phi(z)=\sum_{\sum s_{i} \leq d} \xi(\underline{\mathrm{o}}, \underline{\mathrm{~s}}, \underline{\mathrm{o}})\left(h_{1}-1\right)^{s_{1}} \ldots\left(h_{l}-1\right)^{s_{l}} \\
& \theta^{-1}(z) \equiv \sum \xi(\underline{\mathrm{r}}, \underline{\mathrm{~s}}, \underline{\mathrm{t}}) f_{\beta_{1}}^{r_{1}} \ldots f_{\beta_{N}}^{r_{N}} h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} e_{\beta_{1}}^{t_{1}} \ldots e_{\beta_{N}}^{t_{N}} \quad \bmod S_{d-1}(L) \\
& \eta \theta^{-1}(z) \equiv \sum \xi(\underline{\mathrm{o}}, \underline{\mathrm{~s}}, \underline{\mathrm{o}}) h_{1}^{s_{1}} \ldots h_{l}^{s_{l}} \quad \bmod S_{d-1}(H) .
\end{aligned}
$$

Now it is apparent that

$$
\tau \phi(z) \equiv \phi(z) \quad \bmod S_{d-1}(H)
$$

hence

$$
\tau \phi(z) \equiv \eta \theta^{-1}(z) \quad \bmod S_{d-1}(H)
$$

Since $\tau \phi(z)$ and $\eta \theta^{-1}(z)$ both lie in $S(H)^{W}$ they satisfy $\tau \phi(z) \equiv \eta \theta^{-1}(z)$ $\bmod \left(S(H)^{W}\right)_{d-1}$. Thus $\operatorname{gr}(\tau \phi)=\operatorname{gr}\left(\eta \theta^{-1}\right)$.

Now the maps $\theta^{-1}: Z(L) \rightarrow S(L)^{G}$ and $\eta: S(L)^{G} \rightarrow S(H)^{W}$ satisfy

$$
\begin{aligned}
\theta^{-1}\left(Z_{d}(L)\right) & =\left(S(L)^{G}\right)_{d} \\
\eta\left(S(L)^{G}\right)_{d} & =\left(S(H)^{W}\right)_{d} .
\end{aligned}
$$

Thus we have

$$
\eta \theta^{-1}\left(Z_{d}(L)\right)=\left(S(H)^{W}\right)_{d} .
$$

We may now apply Lemma 11.29. The map

$$
\eta \theta^{-1}: Z(L) \rightarrow S(H)^{W}
$$

is bijective and satisfies

$$
\eta \theta^{-1}\left(Z_{d}(L)\right)=\left(S(H)^{W}\right)_{d}
$$

for each $d$. Hence

$$
\operatorname{gr}\left(\eta \theta^{-1}\right): \operatorname{gr} Z(L) \rightarrow \operatorname{gr} S(H)^{W}
$$

is bijective. This is turn implies that

$$
\tau \phi: Z(L) \rightarrow S(H)^{W}
$$

is bijective. Since $\tau \phi$ is known to be a homomorphism of algebras, it must therefore be an algebra isomorphism.

We can deduce from this theorem a necessary and sufficient condition for two central characters $\chi_{\lambda}, \chi_{\mu}$ to be equal.

Theorem 11.31 Let $\lambda, \mu \in H^{*}$. Then $\chi_{\lambda}=\chi_{\mu}$ if and only if $\mu+\rho=w(\lambda+\rho)$ for some $w \in W$.

Proof. Suppose first that $\mu+\rho=w(\lambda+\rho)$. Then, for $z \in Z(L)$, we have

$$
\begin{aligned}
\chi_{\mu}(z) & =\mu(\phi(z))=(w(\lambda+\rho)-\rho)(\phi(z)) \\
& =w(\lambda+\rho)(\tau(\phi(z)))=(\lambda+\rho)\left(w^{-1} \tau(\phi(z))\right) .
\end{aligned}
$$

Now $\tau \phi(z) \in S(H)^{W}$ and so is fixed by $w^{-1}$. Hence

$$
\chi_{\mu}(z)=(\lambda+\rho)(\tau(\phi(z)))=\lambda(\phi(z))=\chi_{\lambda}(z),
$$

by Theorem 11.25 . Hence $\chi_{\mu}=\chi_{\lambda}$.

Suppose conversely that $\mu+\rho \neq w(\lambda+\rho)$ for all $w \in W$. Then the finite sets $W(\lambda+\rho)$ and $W(\mu+\rho)$ do not intersect. Therefore there exists a polynomial function $Q \in P\left(H^{*}\right)$ such that $Q$ takes values 1 on $W(\lambda+\rho)$ and values 0 on $W(\mu+\rho)$. We have

$$
Q \in S(H)=P\left(H^{*}\right)
$$

By replacing $Q$ by $\frac{1}{|W|} \sum_{w \in W} w(Q)$ we may assume $Q$ lies in $S(H)^{W}$.
We now make use of the isomorphism $\tau \phi: Z(L) \rightarrow S(H)^{W}$. There exists $z \in Z(L)$ such that $\tau \phi(z)=Q$. Thus we have

$$
\begin{aligned}
& \chi_{\lambda}(z)=\lambda(\phi(z))=(\lambda+\rho)(\tau \phi(z))=(\lambda+\rho) Q=1 \\
& \chi_{\mu}(z)=\mu(\phi(z))=(\mu+\rho)(\tau \phi(z))=(\mu+\rho) Q=0
\end{aligned}
$$

Hence $\chi_{\lambda} \neq \chi_{\mu}$.
A second deduction from Theorem 11.30 is the following important result.

Theorem 11.32 The centre $Z(L)$ of $\mathfrak{H}(L)$ is isomorphic to the polynomial ring over $\mathbb{C}$ in $l$ variables, where $L$ is semisimple and $l=\operatorname{rank} L$.

Proof. This follows from Theorem 11.30, Corollary 11.22 and Theorem 11.17.

As an example we consider the Lie algebra $L$ of type $A_{1}$. The algebra $L$ has a basis $f, h, e$ with

$$
[h e]=2 e, \quad[h f]=-2 f, \quad[e f]=h
$$

The algebras

$$
S(H)^{W}, P(H)^{W}, S(L)^{G}, P(L)^{G}, Z(L)
$$

are all isomorphic to the polynomial ring over $\mathbb{C}$ in one variable. We find a generator of each of these algebras.

We have $W=\langle s\rangle$ where $s(h)=-h$. Thus $S(H)^{W}$ is the polynomial algebra generated by $h^{2}$.

We now consider the isomorphism $S(L)^{G} \rightarrow S(H)^{W}$ given by projection. The element of $S(L)^{G}$ mapping to $h^{2}$ is homogeneous of degree 2 in $e, h, f$
and has weight 0 . It must therefore have form $h^{2}+\xi f e$ for some $\xi \in \mathbb{C}$. We determine the constant $\xi$. We have

$$
\text { ad } e \cdot h=-2 e, \quad \text { ad } e \cdot f=h, \quad \text { ad } e \cdot e=0
$$

Thus

$$
\begin{aligned}
& (\exp \operatorname{ad} e) h=h-2 e \\
& (\operatorname{expad} e) f=f+h-e \\
& (\operatorname{expad} e) e=e \\
& \begin{aligned}
(\operatorname{expad} e)\left(h^{2}+\xi f e\right) & =(h-2 e)^{2}+\xi(f+h-e) e \\
& =h^{2}+\xi f e+(\xi-4) h e+(4-\xi) e^{2}
\end{aligned}
\end{aligned}
$$

Thus expad $e$ fixes $h^{2}+\xi f e$ if and only if $\xi=4$. Hence $S(L)^{G}$ is the polynomial ring generated by $h^{2}+4 f e$.

Next we consider the Killing isomorphism $L \rightarrow L^{*}$. $L^{*}$ has basis $f^{*}, h^{*}, e^{*}$ dual to $f, h, e$, that is $y^{*}(x)=1$ if $y=x$ and $y^{*}(x)=0$ if $y \neq x$. Now the Killing form satisfies

$$
\langle h, h\rangle=8, \quad\langle f, e\rangle=4, \quad\langle h, f\rangle=0,\langle h, e\rangle=0, \quad\langle e, e\rangle=0, \quad\langle f, f\rangle=0
$$

Thus under the Killing isomorphism $L \rightarrow L^{*}$ we have $e \rightarrow 4 f^{*}, h \rightarrow 8 h^{*}$, $f \rightarrow 4 e^{*}$. This induces a map $S(L) \rightarrow P(L)$ under which $h^{2}+4 f e$ maps to $64\left(h^{* 2}+f^{*} e^{*}\right)$. Thus $P(L)^{G}$ is the polynomial ring generated by $h^{* 2}+f^{*} e^{*}$.

We also have a map $S(L)^{G} \rightarrow Z(G)$ given by symmetrisation. Under this map $h^{2}+4 f e$ is transformed into

$$
h^{2}+2 f e+2 e f=h^{2}+2 h+4 f e
$$

Thus $Z(L)$ is the polynomial ring generated by $h^{2}+2 h+4 f e$. We also note that the element of $Z(L)$ mapping to $h^{2} \in S(H)^{W}$ under the twisted HarishChandra homomorphism is $h^{2}+2 h+1+4 f e$.

Thus we have:

$$
\begin{aligned}
& S(H)^{W}=\mathbb{C}\left[h^{2}\right] \\
& P(H)^{W}=\mathbb{C}\left[h^{* 2}\right] \\
& S(L)^{G}=\mathbb{C}\left[h^{2}+4 f e\right] \\
& P(L)^{G}=\mathbb{C}\left[h^{* 2}+f^{*} e^{*}\right] \\
& Z(L)=\mathbb{C}\left[h^{2}+2 h+4 f e\right]
\end{aligned}
$$

### 11.6 The Casimir element

We now introduce an element of the centre $Z(L)$ of $\mathfrak{H}(L)$ which has useful properties. Let $x_{1}, \ldots, x_{n}$ be a basis of $L$. Since the Killing form of $L$ is non-degenerate by Theorem 4.10 there is a unique dual basis $y_{1}, \ldots, y_{n}$ of $L$ satisfying

$$
\left\langle x_{i}, y_{j}\right\rangle=\delta_{i j}
$$

Let $c \in \mathfrak{U}(L)$ be defined by

$$
c=\sum_{i=1}^{n} x_{i} y_{i}
$$

Proposition 11.33 The element $c$ is independent of the choice of basis $x_{1}, \ldots, x_{n}$ of $L$.

Proof. Suppose $x_{1}^{\prime}, \ldots, x_{n}^{\prime}$ are a second basis of $L$ and $y_{1}^{\prime}, \ldots, y_{n}^{\prime}$ are the dual basis. Let

$$
x_{i}^{\prime}=\sum_{j} \sigma_{i j} x_{j} \quad y_{i}^{\prime}=\sum_{j} \tau_{i j} y_{j}
$$

Then we have

$$
\left\langle x_{i}^{\prime}, y_{j}^{\prime}\right\rangle=\left\langle\sum_{k} \sigma_{i k} x_{k}, \sum_{l} \tau_{j l} y_{l}\right\rangle=\sum_{k, l} \sigma_{i k} \tau_{j l}\left\langle x_{k}, y_{l}\right\rangle=\sum_{k} \sigma_{i k} \tau_{j k}
$$

Hence if $\sigma=\left(\sigma_{i j}\right), \tau=\left(\tau_{i j}\right)$ we have $\sigma \tau^{\mathrm{t}}=I$. We then have

$$
\sum_{i} x_{i}^{\prime} y_{i}^{\prime}=\sum_{i}\left(\sum_{j} \sigma_{i j} x_{j}\right)\left(\sum_{k} \tau_{i k} y_{k}\right)=\sum_{j, k}\left(\sum_{i} \sigma_{i j} \tau_{i k}\right) x_{j} y_{k} .
$$

Now $\sigma^{\mathrm{t}} \tau=I$ so $\sum_{i} \sigma_{i j} \tau_{i k}=\delta_{j k}$. Hence $\sum_{i} x_{i}^{\prime} y_{i}^{\prime}=\sum_{i} x_{i} y_{i}$.
Definition $c$ is called the Casimir element of $\mathfrak{H}(L)$.
Proposition 11.34 c lies in the centre $Z(L)$ of $\mathfrak{H}(L)$.
Proof. It is sufficient to show that $c x=x c$ for all $x \in L$. We have

$$
\begin{aligned}
c x & =\sum_{i} x_{i} y_{i} x=\sum_{i} x_{i}\left(x y_{i}+\left[y_{i} x\right]\right) \\
& =\sum_{i}\left(\left(x x_{i}+\left[x_{i} x\right]\right) y_{i}+x_{i}\left[y_{i} x\right]\right) \\
& =x c+\sum_{i}\left(\left[x_{i} x\right] y_{i}+x_{i}\left[y_{i} x\right]\right) .
\end{aligned}
$$

Let $\left[x_{i} x\right]=\sum_{j} \alpha_{i j} x_{j}$ and $\left[y_{i} x\right]=\sum_{j} \beta_{i j} y_{j}$. Since $\left\langle\left[x_{i} x\right], y_{j}\right\rangle=\left\langle x_{i},\left[x y_{j}\right]\right\rangle$ we have $\alpha_{i j}=-\beta_{j i}$. It follows that

$$
\begin{aligned}
\sum_{i}\left(\left[x_{i} x\right] y_{i}+x_{i}\left[y_{i} x\right]\right) & =\sum_{i} \sum_{j} \alpha_{i j} x_{j} y_{i}+\sum_{i} \sum_{j} \beta_{i j} x_{i} y_{j} \\
& =\sum_{i, j}\left(\alpha_{i j}+\beta_{j i}\right) x_{j} y_{i}=0 .
\end{aligned}
$$

Thus $c x=x c$ and so $c \in Z(L)$.
We now recall from Proposition 4.18 that for each $e_{\alpha} \in L_{\alpha}$ we can find $f_{\alpha} \in L_{-\alpha}$ such that $\left[e_{\alpha} f_{\alpha}\right]=h_{\alpha}^{\prime}$, and that we then have $\left\langle e_{\alpha}, f_{\alpha}\right\rangle=1$. Since the Killing form of $L$ remains non-degenerate on $H$ we may choose a basis $h_{1}^{\prime}, \ldots, h_{l}^{\prime}$ of $H$ and there will be a dual basis $h_{1}^{\prime \prime}, \ldots, h_{l}^{\prime \prime}$ satisfying

$$
\left\langle h_{i}^{\prime}, h_{j}^{\prime \prime}\right\rangle=\delta_{i j} .
$$

Then $h_{1}^{\prime}, \ldots, h_{l}^{\prime}, e_{\alpha}\left(\alpha \in \Phi^{+}\right), f_{\alpha}\left(\alpha \in \Phi^{+}\right)$are a basis of $L$ and its dual basis is

$$
h_{1}^{\prime \prime}, \ldots, h_{l}^{\prime \prime}, \quad f_{\alpha}\left(\alpha \in \Phi^{+}\right), \quad e_{\alpha}\left(\alpha \in \Phi^{+}\right) .
$$

Using this pair of dual bases we have

$$
c=h_{1}^{\prime} h_{1}^{\prime \prime}+\cdots+h_{l}^{\prime} h_{l}^{\prime \prime}+\sum_{\alpha \in \Phi^{+}} e_{\alpha} f_{\alpha}+\sum_{\alpha \in \Phi^{+}} f_{\alpha} e_{\alpha} .
$$

Thus we obtain:
Proposition 11.35 The Casimir element of $Z(L)$ is given by

$$
c=\sum_{i=1}^{l} h_{i}^{\prime} h_{i}^{\prime \prime}+\sum_{\alpha \in \Phi^{+}} h_{\alpha}^{\prime}+2 \sum_{\alpha \in \Phi^{+}} f_{\alpha} e_{\alpha}
$$

where $h_{1}^{\prime}, \ldots, h_{l}^{\prime} ; h_{1}^{\prime \prime}, \ldots, h_{l}^{\prime \prime}$ are any pair of dual bases of $H$.
The properties of the Casimir element will be useful as we explore further the representation theory of $L$.

Proposition 11.36 Let $c \in Z(L)$ be the Casimir element. Then

$$
\chi_{\lambda}(c)=\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle .
$$

Thus $c$ acts on the Verma module $M(\lambda)$ as scalar multiplication by $\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle$.

Proof. We consider the action of $c$ on the highest weight vector $m_{\lambda}$ of $M(\lambda)$. By Proposition 11.35 we have

$$
\begin{aligned}
c m_{\lambda} & =\left(\sum_{i=1}^{l} h_{i}^{\prime} h_{i}^{\prime \prime}+\sum_{\alpha \in \Phi^{+}} h_{\alpha}^{\prime}+2 \sum_{\alpha \in \Phi^{+}} f_{\alpha} e_{\alpha}\right) m_{\lambda} \\
& =\left(\sum_{i=1}^{l} \lambda\left(h_{i}^{\prime}\right) \lambda\left(h_{i}^{\prime \prime}\right)+\sum_{\alpha \in \Phi^{+}} \lambda\left(h_{\alpha}^{\prime}\right)\right) m_{\lambda}
\end{aligned}
$$

Now $\sum_{\alpha \in \Phi^{+}} \lambda\left(h_{\alpha}^{\prime}\right)=\sum_{\alpha \in \Phi^{+}}\langle\lambda, \alpha\rangle=\left\langle\lambda, \sum \alpha_{\alpha \in \Phi^{+}}\right\rangle=2\langle\lambda, \rho\rangle$.
Let $h_{\lambda}^{\prime} \in H$ be the element corresponding to $\lambda \in H^{*}$ under the isomorphism defined by the Killing form. Thus

$$
\begin{aligned}
\lambda\left(h_{i}^{\prime}\right) & =\left\langle h_{\lambda}^{\prime}, h_{i}^{\prime}\right\rangle \\
\lambda\left(h_{i}^{\prime \prime}\right) & =\left\langle h_{\lambda}^{\prime}, h_{i}^{\prime \prime}\right\rangle
\end{aligned}
$$

We express $h_{\lambda}^{\prime}$ in terms of the dual bases $h_{1}^{\prime}, \ldots, h_{l}^{\prime}$ and $h_{1}^{\prime \prime}, \ldots, h_{l}^{\prime \prime}$ of $H$. Let

$$
\begin{aligned}
& h_{\lambda}^{\prime}=a_{1} h_{1}^{\prime}+\cdots+a_{l} h_{l}^{\prime} \\
& h_{\lambda}^{\prime}=b_{1} h_{1}^{\prime \prime}+\cdots+b_{l} h_{l}^{\prime \prime}
\end{aligned}
$$

Since $\left\langle h_{i}^{\prime}, h_{j}^{\prime \prime}\right\rangle=\delta_{i j}$ we have

$$
\begin{aligned}
& \left\langle h_{\lambda}^{\prime}, h_{\lambda}^{\prime}\right\rangle=a_{1} b_{1}+\cdots+a_{l} b_{l} \\
& \left\langle h_{\lambda}^{\prime}, h_{i}^{\prime}\right\rangle=b_{i} \quad\left\langle h_{\lambda}^{\prime}, h_{i}^{\prime \prime}\right\rangle=a_{i}
\end{aligned}
$$

It follows that

$$
\sum_{i=1}^{l} \lambda\left(h_{i}^{\prime}\right) \lambda\left(h_{i}^{\prime \prime}\right)=\sum_{i=1}^{l}\left\langle h_{\lambda}^{\prime}, h_{i}^{\prime}\right\rangle\left\langle h_{\lambda}^{\prime}, h_{i}^{\prime \prime}\right\rangle=\left\langle h_{\lambda}^{\prime}, h_{\lambda}^{\prime}\right\rangle=\langle\lambda, \lambda\rangle .
$$

Hence

$$
\begin{aligned}
c m_{\lambda} & =(\langle\lambda, \lambda\rangle+2\langle\lambda, \rho\rangle) m_{\lambda} \\
& =(\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle) m_{\lambda} .
\end{aligned}
$$

Thus the value of the central character $\chi_{\lambda}$ at $c$ is given by

$$
\chi_{\lambda}(c)=\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle .
$$

## 12

## Character and dimension formulae

### 12.1 Characters of $L$-modules

Let $V$ be an $L$-module where $L$ is semisimple. We say that $V$ admits a character if $V$ is the direct sum of its weight spaces and each weight space of $V$ is finite dimensional. Thus we have

$$
V=\bigoplus_{\lambda \in H^{*}} V_{\lambda} \quad \operatorname{dim} V_{\lambda} \quad \text { finite }
$$

where $V_{\lambda}=\{v \in V ; h v=\lambda(h) v$ for all $h \in H\}$. The character of $V$ is then the function ch $V: H^{*} \rightarrow \mathbb{Z}$ given by

$$
(\operatorname{ch} V)(\lambda)=\operatorname{dim} V_{\lambda} .
$$

We see that if $V$ admits a character then the structure of $V$ as an $H$-module is determined by ch $V$.

In this chapter we shall obtain formulae for the characters of the Verma modules $M(\lambda)$ for $\lambda \in H^{*}$ and for the finite dimensional irreducible modules $L(\lambda)$ for $\lambda \in X^{+}$.

We first identify a certain ring of functions $H^{*} \rightarrow \mathbb{Z}$ in which it will be convenient to work. Given a function $f: H^{*} \rightarrow \mathbb{Z}$ we define Supp $f$, the support of $f$, to be the set of $\lambda \in H^{*}$ for which $f(\lambda) \neq 0$. For example the support of the function ch $M(\lambda)$ is the set of all $\mu \in H^{*}$ which have form

$$
\mu=\lambda-n_{1} \alpha_{1}-\cdots-n_{l} \alpha_{l} \quad n_{i} \in \mathbb{Z}, \quad n_{i} \geq 0
$$

This follows from Theorem 10.7. We define

$$
S(\lambda)=\operatorname{Supp}(\operatorname{ch} M(\lambda))
$$

Definition $\mathfrak{R}$ denotes the set of all functions $f: H^{*} \rightarrow \mathbb{Z}$ such that there exists a finite set $\lambda_{1}, \ldots, \lambda_{k} \in H^{*}$ with

$$
\operatorname{Supp} f \subset S\left(\lambda_{1}\right) \cup \cdots \cup S\left(\lambda_{k}\right)
$$

It is clear that $\operatorname{ch} M(\lambda)$ for $\lambda \in H^{*}$ and $\operatorname{ch} L(\lambda)$ for $\lambda \in X^{+}$lie in $\Re$. It is also clear that if $f, g \in \mathfrak{R}$ then $f+g \in \mathfrak{R}$, since

$$
\operatorname{Supp}(f+g) \subset \operatorname{Supp} f \cup \operatorname{Supp} g .
$$

Thus $\mathfrak{R}$ is an additive group. We can also define a product on $\mathfrak{R}$ which makes it into a ring. Given $f, g \in \mathfrak{R}$ we define $f g: H^{*} \rightarrow \mathbb{Z}$ by

$$
(f g)(\lambda)=\sum_{\substack{\mu, \nu \in H^{*} \\ \mu+\nu=\lambda}} f(\mu) g(\nu)
$$

We note that the sum is finite, so that $f g$ is well defined. For we may assume $\mu \in \operatorname{Supp} f$ and $\nu \in \operatorname{Supp} g$. Suppose

$$
\begin{aligned}
& \operatorname{Supp} f \subset S\left(\mu_{1}\right) \cup \cdots \cup S\left(\mu_{h}\right) \\
& \operatorname{Supp} g \subset S\left(\nu_{1}\right) \cup \cdots \cup S\left(\nu_{h}\right)
\end{aligned}
$$

If $\mu \in S\left(\mu_{i}\right)$ and $\nu \in S\left(\nu_{j}\right)$ we have

$$
\begin{aligned}
\mu & =\mu_{i}-m_{1} \alpha_{1}-\cdots-m_{l} \alpha_{l} & m_{k} \in \mathbb{Z}, & m_{k} \geq 0 \\
\nu & =\nu_{j}-n_{1} \alpha_{1}-\cdots-n_{l} \alpha_{l} & n_{k} \in \mathbb{Z}, & n_{k} \geq 0
\end{aligned}
$$

Since $\mu+\nu=\lambda$ we have

$$
\lambda=\left(\mu_{i}+\nu_{j}\right)-r_{1} \alpha_{1}-\cdots-r_{l} \alpha_{l} \quad r_{k} \in \mathbb{Z}, r_{k} \geq 0
$$

where $r_{k}=m_{k}+n_{k}$. However, given $i, j$ and $\lambda$ the non-negative integers $m_{k}, n_{k}$ with $m_{k}+n_{k}=r_{k}$ can be chosen in only finitely many ways, thus our sum is finite. Also we see that $\operatorname{Supp}(f g) \subset \bigcup_{i, j} S\left(\mu_{i}+\nu_{j}\right)$, hence $f g \in \mathfrak{R}$. It is also readily checked that $(f g) h=f(g h)$, thus $\mathfrak{R}$ becomes a ring.

For each $\lambda \in H^{*}$ we define $e_{\lambda}: H^{*} \rightarrow \mathbb{Z}$ by $e_{\lambda}(\lambda)=1, e_{\lambda}(\mu)=0$ if $\mu \neq \lambda$. Thus $e_{\lambda}$ is the characteristic function of $\lambda$. All such characteristic functions lie in $\mathfrak{R}$. In fact if $f$ is any function in $\mathfrak{R}$ it is convenient to write

$$
f=\sum_{\lambda \in H^{*}} f(\lambda) e_{\lambda}
$$

even though the sum may be infinite.
We note that $e_{\lambda} e_{\mu}=e_{\lambda+\mu}$.

Lemma 12.1 Suppose that the L-module $V$ admits a character and let $U$ be a submodule of $V$. Then both $U$ and $V / U$ admit a character, and

$$
\operatorname{ch} U+\operatorname{ch} \frac{V}{U}=\operatorname{ch} V
$$

Proof. We have $V=\bigoplus_{\lambda} V_{\lambda}$. Also $U_{\lambda}=U \cap V_{\lambda}$. Thus the sum $\sum_{\lambda} U_{\lambda}$ is direct. Moreover $U=\sum_{\lambda} U_{\lambda}$ since if $u \in U$ and $u=\sum u_{\lambda}$ with $u_{\lambda} \in V_{\lambda}$ then $u_{\lambda} \in U$, as in the proof of Theorem 10.9. Hence we have

$$
U=\bigoplus_{\lambda} U_{\lambda}
$$

with $U_{\lambda} \subset V_{\lambda}$, so $U$ admits a character. We also have

$$
V / U=\bigoplus_{\lambda}\left(V_{\lambda} / U_{\lambda}\right)
$$

and $V_{\lambda} / U_{\lambda}$ can be identified with the $\lambda$-weight space $(V / U)_{\lambda}$. Thus $V / U$ admits a character. Finally we have

$$
(\operatorname{ch} U)(\lambda)+(\operatorname{ch}(V / U))(\lambda)=\operatorname{dim} U_{\lambda}+\operatorname{dim}\left(V_{\lambda} / U_{\lambda}\right)=\operatorname{dim} V_{\lambda} .
$$

Thus ch $U+\operatorname{ch}(V / U)=\operatorname{ch} V$.
Lemma 12.2 Suppose $V_{1}, V_{2}$ are L-modules which both admit characters such that ch $V_{1}$ and ch $V_{2}$ lie in $\mathfrak{R}$. Then $V_{1} \otimes V_{2}$ admits a character and $\operatorname{ch}\left(V_{1} \otimes V_{2}\right)=\operatorname{ch} V_{1} \operatorname{ch} V_{2}$.

Proof. Since $V_{1}, V_{2}$ admit characters we have $V_{1}=\bigoplus_{\lambda}\left(V_{1}\right)_{\lambda}$ and $V_{2}=\bigoplus_{\mu}\left(V_{2}\right)_{\mu}$. Hence

$$
V_{1} \otimes V_{2}=\bigoplus_{\lambda, \mu}\left(\left(V_{1}\right)_{\lambda} \otimes\left(V_{2}\right)_{\mu}\right) .
$$

$V_{1} \otimes V_{2}$ may be made into an $L$-module by means of the action

$$
x\left(v_{1} \otimes v_{2}\right)=x v_{1} \otimes v_{2}+v_{1} \otimes x v_{2}
$$

extended by linearity. In particular, if $x \in H, v_{1} \in\left(V_{1}\right)_{\lambda}$ and $v_{2} \in\left(V_{2}\right)_{\mu}$ we have

$$
x\left(v_{1} \otimes v_{2}\right)=(\lambda(x)+\mu(x)) v_{1} \otimes v_{2} .
$$

Thus $\left(V_{1}\right)_{\lambda} \otimes\left(V_{2}\right)_{\mu} \subset\left(V_{1} \otimes V_{2}\right)_{\lambda+\mu}$. It follows that

$$
V_{1} \otimes V_{2}=\oplus\left(V_{1} \otimes V_{2}\right)_{\nu}
$$

where $\left(V_{1} \otimes V_{2}\right)_{\nu}=\sum_{\substack{\lambda, \mu \\ \lambda+\mu=\nu}}\left(\left(V_{1}\right)_{\lambda} \otimes\left(V_{2}\right)_{\mu}\right)$. Thus $V_{1} \otimes V_{2}$ admits a character. Moreover we have

$$
\begin{aligned}
\left(\operatorname{ch}\left(V_{1} \otimes V_{2}\right)\right)(\nu) & =\operatorname{dim}\left(V_{1} \otimes V_{2}\right)_{\nu}=\sum_{\substack{\lambda, \mu \\
\lambda+\mu=\nu}} \operatorname{dim}\left(V_{1}\right)_{\lambda} \operatorname{dim}\left(V_{2}\right)_{\mu} \\
& =\sum_{\substack{\lambda, \mu \\
\lambda+\mu=\nu}}\left(\operatorname{ch} V_{1}\right)(\lambda)\left(\operatorname{ch} V_{2}\right)(\mu)=\left(\operatorname{ch} V_{1} \operatorname{ch} V_{2}\right)(\nu)
\end{aligned}
$$

Thus ch $\left(V_{1} \otimes V_{2}\right)=\operatorname{ch} V_{1} \operatorname{ch} V_{2}$ as required.

### 12.2 Characters of Verma modules

We now consider the character of the Verma module $M(\lambda)$ where $\lambda \in H^{*}$. We recall from Theorem 10.7 that

$$
(\operatorname{ch} M(\lambda))(\mu)=\mathfrak{B}(\lambda-\mu)
$$

where $\mathfrak{P}(\lambda-\mu)$ is the number of ways of expressing $\lambda-\mu$ as a sum of positive roots. Thus we have

$$
\begin{aligned}
\operatorname{ch} M(\lambda) & =\sum_{\mu \in H^{*}} \mathfrak{B}(\lambda-\mu) e_{\mu}=\sum_{\nu \in H^{*}} \mathfrak{P}(\nu) e_{\lambda-\nu} \\
& =\sum_{\nu \in H^{*}} \mathfrak{P}(\nu) e_{\lambda} e_{-\nu}=e_{\lambda} \sum_{\nu \in H^{*}} \mathfrak{P}(\nu) e_{-\nu}
\end{aligned}
$$

We write $\Gamma=\sum_{\nu \in H^{*}} \mathfrak{B}(\nu) e_{-\nu}$. We have $\Gamma \in \mathfrak{R}$ since $\operatorname{Supp} \Gamma \subset S(0)$. Then we have

$$
\operatorname{ch} M(\lambda)=e_{\lambda} \Gamma
$$

Lemma $12.3 \Gamma$ has an inverse in the ring $\mathfrak{R}$ given by

$$
\Gamma^{-1}=\prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)
$$

Proof. Let $\Phi^{+}=\left\{\beta_{1}, \ldots, \beta_{N}\right\}$. Then $\mathfrak{X}(\nu) \neq 0$ if and only if there exist nonnegative integers $r_{1}, \ldots, r_{N}$ such that $\nu=r_{1} \beta_{1}+\cdots+r_{N} \beta_{N}$. In fact $\mathfrak{B}(\nu)$ is the number of such sets $\left(r_{1}, \ldots, r_{N}\right)$. Thus we have

$$
\begin{aligned}
\Gamma & =\sum_{\nu} \mathfrak{B}(\nu) e_{-\nu}=\sum_{r_{1}, \ldots, r_{N} \geq 0} e_{-r_{1} \beta_{1}-\ldots-r_{N} \beta_{N}} \\
& =\sum_{r_{1}, \ldots, r_{N} \geq 0} e_{-\beta_{1}}^{r_{1}} \ldots e_{-\beta_{N}}^{r_{N}}=\prod_{i=1}^{N}\left(\sum_{r_{i} \geq 0} e_{-\beta_{i}}^{r_{i}}\right) .
\end{aligned}
$$

This factorisation of $\Gamma$ in $\mathfrak{R}$ gives us the required result. For the element

$$
1+e_{-\beta_{i}}+e_{-\beta_{i}}^{2}+\cdots \quad \text { of } \mathfrak{R}
$$

has an inverse $1-e_{-\beta_{i}} \in \mathfrak{R}$. Thus $\Gamma$ has an inverse

$$
\Gamma^{-1}=\prod_{i=1}^{N}\left(1-e_{-\beta_{i}}\right)=\prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)
$$

This gives us a useful formula for the character of the Verma module $M(\lambda)$.

Proposition $12.4 \operatorname{ch} M(\lambda)=e_{\frac{\lambda+\rho}{\Delta}}$ where $\Delta=e_{\rho} \prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)$.
Proof. We have

$$
\operatorname{ch} M(\lambda)=e_{\lambda} \Gamma=\frac{e_{\lambda} e_{\rho}}{\Delta}=\frac{e_{\lambda+\rho}}{\Delta}
$$

by Lemma 12.3.
The denominator $\Delta$ is an element of $\mathfrak{R}$ which can be expressed in a number of alternative ways.

We recall that $\rho \in X$ was defined by

$$
\rho=\omega_{1}+\cdots+\omega_{l}
$$

i.e. $\rho$ is the sum of the fundamental weights. This element can also be expressed simply in terms of the roots.

Proposition 12.5 $\rho=\frac{1}{2} \sum_{\alpha \in \Phi^{+}} \alpha$. Thus $\rho$ is one half the sum of the positive roots.

Proof. Let $\rho^{\prime}=\frac{1}{2} \sum_{\alpha \in \Phi^{+}} \alpha$. We can express $\rho^{\prime}$ as a linear combination of the fundamental weights. Let

$$
\rho^{\prime}=\sum_{i=1}^{l} c_{i} \omega_{i} \quad \text { with } c_{i} \in \mathbb{Q}
$$

Now the fundamental reflection $s_{i} \in W$ transforms $\alpha_{i}$ to $-\alpha_{i}$ and transforms every other positive root to a positive root, by Lemma 5.9. Thus we have

$$
s_{i}\left(\rho^{\prime}\right)=\rho^{\prime}-\alpha_{i}
$$

On the other hand we have

$$
2 \frac{\left\langle\alpha_{j}, \omega_{i}\right\rangle}{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}=\delta_{i j}
$$

by Proposition 10.18. This shows that $s_{j}\left(\omega_{i}\right)=\omega_{i}$ if $i \neq j$ and $s_{j}\left(\omega_{j}\right)=\omega_{j}-\alpha_{j}$. Thus we have

$$
s_{i}\left(\rho^{\prime}\right)=\rho^{\prime}-c_{i} \alpha_{i}
$$

Comparing this with the above formula for $s_{i}\left(\rho^{\prime}\right)$ we deduce that $c_{i}=1$. Hence $\rho^{\prime}=\rho$ as required.

Corollary $12.6 \Delta=e_{-\rho} \prod_{\alpha \in \Phi^{+}}\left(e_{\alpha}-1\right)$

Proof. We have

$$
\begin{aligned}
\Delta & =e_{\rho} \prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)=e_{\rho} \prod_{\alpha \in \Phi^{+}} e_{-\alpha}\left(e_{\alpha}-1\right)=e_{\rho}\left(\prod_{\alpha \in \Phi^{+}} e_{-\alpha}\right) \prod_{\alpha \in \Phi^{+}}\left(e_{\alpha}-1\right) \\
& =e_{\rho} e_{-2 \rho} \prod_{\alpha \in \Phi^{+}}\left(e_{\alpha}-1\right)=e_{-\rho} \prod_{\alpha \in \Phi^{+}}\left(e_{\alpha}-1\right) .
\end{aligned}
$$

There is a further useful expression for the denominator $\Delta$. Before proving it we shall need some information about the geometry of the action of the Weyl group $W$ on the Euclidean space $V=H_{\mathbb{R}}^{*}$.

### 12.3 Chambers and roots

We recall that the Weyl group is a finite group of isometries of the Euclidean space $V$ generated by the reflections $s_{\alpha}$ for $\alpha \in \Phi$. We have

$$
s_{\alpha}(v)=v-2 \frac{\langle\alpha, v\rangle}{\langle\alpha, \alpha\rangle} \alpha \quad v \in V
$$

Let

$$
\begin{aligned}
L_{\alpha} & =\left\{v \in V ; s_{\alpha}(v)=v\right\} \\
& =\{v \in V ;\langle\alpha, v\rangle=0\} .
\end{aligned}
$$

$L_{\alpha}$ is the reflecting hyperplane orthogonal to the root $\alpha$. We consider the complement

$$
V-\cup_{\alpha \in \Phi} L_{\alpha}
$$

of the set of reflecting hyperplanes. This is an open subset of $V$. The connected components of this set are called the chambers of $V$. Two points of
$V-\bigcup_{\alpha \in \Phi} L_{\alpha}$ lie in the same chamber if and only if they lie on the same side of each reflecting hyperplane.

Let $C$ be a chamber in $V$ and $\delta(C)$ be the boundary of $C$. Then the hyperplanes $L_{\alpha}$ such that $L_{\alpha} \cap \delta(C)$ is not contained in any proper subspace of $L_{\alpha}$ are called the bounding hyperplanes, or walls, of $C$.

Now let $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$ be a fundamental system of roots. Then the set

$$
C=\left\{v \in V ;\left\langle\alpha_{i}, v\right\rangle>0 \quad \text { for } i=1, \ldots, l\right\}
$$

is a chamber of $V$. For if $\alpha$ is any positive root we have $\langle\alpha, v\rangle>0$ for all $v \in C$. Thus all elements of $C$ lie on the same side of each reflecting hyperplane $L_{\alpha}$. Thus $C$ lies in $V-\bigcup_{\alpha \in \Phi} L_{\alpha}$ and $C$ is connected. Moreover any subset of $V-\bigcup_{\alpha \in \Phi} L_{\alpha}$ larger than $C$ would contain an element $v$ with $\left\langle\alpha_{i}, v\right\rangle<0$ for some $i$, and so would be disconnected. $C$ is called the fundamental chamber corresponding to the fundamental system $\Pi$. The bounding hyperplanes of $C$ are $L_{\alpha_{1}}, \ldots, L_{\alpha_{l}}$. For $L_{\alpha_{i}} \cap \delta(C)$ consists of all $v \in V$ such that $\left\langle\alpha_{i}, v\right\rangle=0$ but $\left\langle\alpha_{j}, v\right\rangle \geq 0$ for $j \neq i$. Since $\alpha_{1}, \ldots, \alpha_{l}$ are linearly independent $L_{\alpha_{i}} \cap \delta(C)$ is not contained in any proper subspace of $L_{\alpha_{i}}$. On the other hand let $\alpha$ be a positive root which is not fundamental. Then $\alpha=\sum_{i=1} n_{i} \alpha_{i}$ with each $n_{i} \geq 0$ and at least two $n_{i}>0$. If $v \in L_{\alpha} \cap \delta(C)$ then

$$
\sum n_{i}\left\langle\alpha_{i}, v\right\rangle=0
$$

and so $\left\langle\alpha_{i}, v\right\rangle=0$ whenever $n_{i}>0$. Thus $L_{\alpha} \cap \delta(C)$ lies in a proper subspace of $L_{\alpha}$. Hence the bounding hyperplanes of $C$ are $L_{\alpha_{1}}, \ldots, L_{\alpha_{l}}$. In fact the set $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$ of fundamental roots may be characterised as the roots orthogonal to the bounding hyperplanes of $C$ which point into $C$, that is such that $\alpha_{i}$ lies on the same side of $L_{\alpha_{i}}$ as $C$.

Now the Weyl group acts on $V$ in a way which permutes the roots. It therefore permutes the reflecting hyperplanes $L_{\alpha}$, and so acts on $V-\bigcup_{\alpha \in \Phi} L_{\alpha}$. Since $W$ is a group of isometries of $V, W$ permutes the connected components of $V-\bigcup_{\alpha \in \Phi} L_{\alpha}$. Thus the Weyl group $W$ acts on the set of chambers of $V$.

Proposition 12.7 (i) Given any two chambers $C, C^{\prime}$ of $V$ there is a unique element $w \in W$ such that $w(C)=C^{\prime}$.
(ii) The number of chambers of $V$ is equal to the order of the Weyl group.
(iii) If $C$ is a chamber in $V$ its closure $\bar{C}$ contains just one element from each $W$-orbit on $V$.

Proof. Let $\Pi$ be a fundamental system of roots and $C$ be the chamber defined by $v \in C$ if and only if $\left\langle\alpha_{i}, v\right\rangle>0$ for $i=1, \ldots, l$. Let $C^{\prime}$ be any chamber and let $v \in C^{\prime}$. We recall from Section 5.1 that $\Pi$ is associated with a total
ordering $>$ on $V$. We consider the set of transforms $w(v)$ for $w \in W$ and let $v^{\prime}$ be the one which is greatest in the above total ordering. Then we have

$$
s_{\alpha_{i}}\left(v^{\prime}\right)=v^{\prime}-2 \frac{\left\langle\alpha_{i}, v^{\prime}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} \alpha_{i} \quad \alpha_{i} \in \Pi
$$

and since $s_{\alpha_{i}}\left(v^{\prime}\right) \leq v^{\prime}$ we must have $\left\langle\alpha_{i}, v^{\prime}\right\rangle \geq 0$. This holds for all $i=1, \ldots, l$, thus $v^{\prime} \in C$. Now let $v^{\prime}=w(v)$. Since $v \in C^{\prime}$ we have $v^{\prime} \in w\left(C^{\prime}\right)$. Thus $w\left(C^{\prime}\right)$ is a chamber which intersects $\bar{C}$. However, the only chamber intersecting $\bar{C}$ is $C$. Thus $w\left(C^{\prime}\right)=C$. Hence any chamber $C^{\prime}$ is in the same $W$-orbit as $C$. Thus $W$ acts transitively on the set of chambers. It follows that any chamber is associated to some fundamental system of roots in the manner described above.

Now suppose $w(C)=C$. Then we have $w(\Pi)=\Pi$ where $\Pi$ is the fundamental system determined by $C$, i.e. the set of roots orthogonal to the walls of $C$ and pointing into $C$. It follows that $w\left(\Phi^{+}\right)=\Phi^{+}$, so $w$ makes every positive root positive. Hence $n(w)=0$. It follows from Corollary 5.16 that $l(w)=0$, i.e. $w=1$. Thus $W$ acts simply transitively on the set of chambers.

It is a consequence of this that the number of chambers of $V$ is equal to $|W|$.

We now consider the closure $\bar{C}$ of a chamber $C$. Since each vector lies in the closure of some chamber and $W$ acts transitively on the chambers each orbit of $W$ on $V$ intersects $\bar{C}$. We must also show that if $v_{1}, v_{2} \in \bar{C}$ and $w\left(v_{1}\right)=v_{2}$ then $v_{1}=v_{2}$. We prove this by induction on $l(w)$. It is clear when $l(w)=0$, i.e. $w=1$. Thus we assume $l(w)>0$. Then $n(w)>0$ so there exists $\alpha_{i} \in \Pi$ with $w\left(\alpha_{i}\right)<0$. Thus

$$
0 \leq\left\langle v_{1}, \alpha_{i}\right\rangle=\left\langle v_{2}, w\left(\alpha_{i}\right)\right\rangle \leq 0
$$

Hence $\left\langle v_{1}, \alpha_{i}\right\rangle=0$ and $s_{\alpha_{i}}\left(v_{1}\right)=v_{1}$. But now $w s_{\alpha_{i}}\left(v_{1}\right)=v_{2}$. The only positive root made negative by $s_{\alpha_{i}}$ is $\alpha_{i}$. Thus the positive roots made negative by $w$ and $w s_{\alpha_{i}}$ are the same, apart from $\alpha_{i}$, which is made negative by $w$ and positive by $w s_{\alpha_{i}}$. Thus

$$
n(w)=n\left(w s_{\alpha_{i}}\right)+1
$$

and so

$$
l\left(w s_{\alpha_{i}}\right)=l(w)-1
$$

by Corollary 5.16 . We can then deduce that $v_{1}=v_{2}$ by induction, as required.

We shall now suppose that $\Pi$ is a fixed fundamental system of roots and $C$ is the corresponding fundamental chamber.

Proposition 12.8 (i) $v \in C$ if and only if $v=\sum_{i=1}^{l} n_{i} \omega_{i}$ with $n_{i}>0$ for all $i$.
(ii) $v \in \bar{C}$ if and only if $v=\sum_{i=1}^{l} n_{i} \omega_{i}$ with $n_{i} \geq 0$ for all $i$.

Proof. Since $\omega_{1}, \ldots, \omega_{l}$ are a basis of $V$ we can write $v=\sum n_{i} \omega_{i}$ for each $v \in V$. Now $v \in C$ if and only if $\left\langle\alpha_{i}, v\right\rangle>0$ for $i=1, \ldots, l$. We recall from the definition of the fundamental weights $\omega_{1}, \ldots, \omega_{l}$ that

$$
\begin{aligned}
& \left\langle\alpha_{i}, \omega_{j}\right\rangle=0 \quad \text { if } i \neq j \\
& \left\langle\alpha_{i}, \omega_{i}\right\rangle=2\left\langle\alpha_{i}, \alpha_{i}\right\rangle
\end{aligned}
$$

Thus we have $\left\langle\alpha_{i}, v\right\rangle=2 n_{i}\left\langle\alpha_{i}, \alpha_{i}\right\rangle$. In particular $\left\langle\alpha_{i}, v\right\rangle>0$ if and only if $n_{i}>0$. Similarly $\left\langle\alpha_{i}, v\right\rangle \geq 0$ if and only if $n_{i} \geq 0$. The required result follows.

We show in Figures 12.1, 12.2 and 12.3 the chambers for the 2-dimensional root systems $A_{2}, B_{2}$ and $G_{2}$.


Figure 12.1 Two-dimensional root system type $A_{2}$


Figure 12.2 Two-dimensional root system type $B_{2}$


Figure 12.3 Two-dimensional root system type $G_{2}$

Proposition 12.9 Suppose $\Phi$ is the root system of a simple Lie algebra and let $C$ be the fundamental chamber.
(i) Suppose all roots in $\Phi$ have the same length. Then there exists a unique root $\theta_{l}=\sum_{i=1}^{l} a_{i} \alpha_{i}$ in $\bar{C}$. This root satisfies the condition that for any root $\alpha=\sum_{i=1}^{l} k_{i} \alpha_{i}$ we have $k_{i} \leq a_{i}$.
(ii) Now suppose there are two root lengths. Then there are just two roots

$$
\theta_{\mathrm{l}}=\sum_{i=1}^{l} a_{i} \alpha_{i}, \quad \theta_{\mathrm{s}}=\sum_{i=1}^{l} c_{i} \alpha_{i}
$$

in $\bar{C} . \theta_{1}$ is a long root and $\theta_{\mathrm{s}}$ is a short root. $\theta_{1}$ satisfies the condition that for any root $\alpha=\sum_{i=1}^{l} k_{i} \alpha_{i}$ we have $k_{i} \leq a_{i}$. (In particular $c_{i} \leq a_{i}$.) $\theta_{l}$ is called the highest root and $\theta_{\mathrm{s}}$ the highest short root.

Proof. By Proposition $12.7 \bar{C}$ contains just one root in each $W$-orbit on $\Phi$. Now two roots lie in the same $W$-orbit if and only if they have the same length. For roots in the same orbit obviously have the same length; but any root is in the same orbit as a fundamental root, and any two fundamental roots of the same length can be joined in the Dynkin diagram by a sequence of fundamental roots all of this length. Two fundamental roots of the same length joined in the Dynkin diagram obviously lie in the same $W$-orbit.

Thus in case (i) $\bar{C}$ contains a unique root $\theta_{1}$ and in case (ii) $\bar{C}$ contains one long root $\theta_{1}$ and one short root $\theta_{\mathrm{s}}$.

We now introduce a partial order $\triangleright$ on the set $\Phi^{+}$of positive roots. Given

$$
\alpha=\sum_{i=1}^{l} m_{i} \alpha_{i}, \quad \beta=\sum_{i=1}^{l} n_{i} \alpha_{i}
$$

in $\Phi^{+}$we write $\alpha \triangleright \beta$ if $m_{i} \geq n_{i}$ for each $i$. We consider maximal elements of $\Phi^{+}$with respect to this partial order. Let $\alpha$ be maximal. Then $\left\langle\alpha, \alpha_{i}\right\rangle \geq 0$ for each $i$, as otherwise $\alpha+\alpha_{i}$ would be a root higher than $\alpha$. We also have $\left\langle\alpha, \alpha_{i}\right\rangle>0$ for some $i$. Let $\alpha=\sum_{i=1}^{l} m_{i} \alpha_{i}$. We show that each $m_{i}>0$. Suppose this is not so. Then there exist $i, i^{\prime}$ with $m_{i} \neq 0, m_{i^{\prime}}=0$ and $\left\langle\alpha_{i}, \alpha_{i^{\prime}}\right\rangle<0$. But then $\left\langle\alpha, \alpha_{i^{\prime}}\right\rangle=\sum_{j=1}^{l} m_{j}\left\langle\alpha_{j}, \alpha_{i^{\prime}}\right\rangle<0$, a contradiction. Hence each $m_{i}>0$.

We now show that $\alpha$ is the unique maximal element of $\Phi^{+}$with respect to $\triangleright$. Suppose if possible that $\beta$ is also maximal and $\beta \neq \alpha$. Then $\alpha+\beta \notin \Phi$. Also $\alpha-\beta \notin \Phi$, as $\alpha-\beta \notin \Phi$ would imply $\alpha \triangleright \beta$ or $\beta \triangleright \alpha$. Hence $\langle\alpha, \beta\rangle=0$ by Proposition 4.22. But

$$
\langle\alpha, \beta\rangle=\sum_{i=1}^{l} m_{i}\left\langle\alpha_{i}, \beta\right\rangle>0
$$

since each $m_{i}>0$, each $\left\langle\alpha_{i}, \beta\right\rangle \geq 0$, and some $\left\langle\alpha_{i}, \beta\right\rangle>0$. Thus we have a contradiction. Hence $\alpha$ is the unique maximal element of $\Phi^{+}$with respect to $\triangleright$.
Now $\alpha \in \bar{C}$ since $\left\langle\alpha, \alpha_{i}\right\rangle \geq 0$ for each $i$. Thus $\alpha=\theta_{l}$ or $\theta_{\mathrm{s}}$. We wish to show $\alpha=\theta_{l}$. To do so we show that if $\alpha^{\prime} \in \Phi \cap \bar{C}$ then $\left\langle\alpha^{\prime}, \alpha^{\prime}\right\rangle \leq\langle\alpha, \alpha\rangle$. By the maximality of $\alpha, \alpha-\alpha^{\prime}$ is a non-negative combination of $\alpha_{1}, \ldots, \alpha_{l}$ and so $\left\langle\alpha-\alpha^{\prime}, x\right\rangle \geq 0$ for all $x \in \bar{C}$. In particular we have $\left\langle\alpha-\alpha^{\prime}, \alpha\right\rangle \geq 0$ and $\left\langle\alpha-\alpha^{\prime}, \alpha^{\prime}\right\rangle \geq 0$. Hence

$$
\langle\alpha, \alpha\rangle \geq\left\langle\alpha, \alpha^{\prime}\right\rangle \geq\left\langle\alpha^{\prime}, \alpha^{\prime}\right\rangle
$$

It follows that $\alpha=\theta_{l}$. Thus $\theta_{l}$ is the unique maximal element of $\Phi^{+}$with respect to $\triangleright$ and the result is proved.

Definition The number $h=1+\mathrm{ht} \theta_{l}$ is called the Coxeter number of L. It is known to be equal to the order of the element $s_{1} s_{2} \ldots s_{l} \in W$, and also to $|\Phi| /|\Pi|$. (See, for example, Bourbaki, Groupes et algèbres de Lie, Chapters 4, 5, 6.)

In order to prove Weyl's denominator formula we shall need some properties of the transforms $w(\rho)$.

Proposition 12.10 (i) $w(\rho)=\rho-\sum_{\alpha \in \Omega} \alpha$ for some subset $\Omega$ of $\Phi^{+}$.
(ii) Given any subset $\Omega$ of $\Phi^{+}$the vector $\rho-\sum_{\alpha \in \Omega} \alpha$ either lies in one of the reflecting hyperplanes $L_{\alpha}$ or has the form $w(\rho)$ for some $w \in W$.
(iii) If $\rho-\sum_{\alpha \in \Omega} \alpha$ lies in the fundamental chamber then $\Omega$ is empty.

Proof. We know from Proposition 12.5 that $\rho=\frac{1}{2} \sum_{\alpha \in \Phi^{+}} \alpha$. Let $w \in W$. Then $w$ permutes the roots and so

$$
w(\rho)=\frac{1}{2} \sum_{\alpha \in \Phi^{+}}( \pm \alpha)=\rho-\sum_{\alpha \in \Omega} \alpha
$$

where $\Omega$ is the set of positive roots made negative by $w^{-1}$.
Now suppose $\Omega$ is any subset of $\Phi^{+}$. Suppose $\rho-\sum_{\alpha \in \Omega} \alpha$ lies in the fundamental chamber. We write $v=\sum_{\alpha \in \Omega} \alpha$. Then $(\rho-v)\left(h_{i}\right)>0$ for each $i=1, \ldots, l$. Moreover $(\rho-v)\left(h_{i}\right) \in \mathbb{Z}$ and $\rho\left(h_{i}\right)=1$, since $\omega_{i}\left(h_{i}\right)=1$ and $\omega_{j}\left(h_{i}\right)=0$ if $j \neq i$. It follows that $v\left(h_{i}\right) \leq 0$, that is that $\left\langle v, \alpha_{i}\right\rangle \leq 0$ for $i=$ $1, \ldots, l$. However, $v$ is a sum of positive roots so has form $v=\sum_{i=1}^{l} n_{i} \alpha_{i}$ where $n_{i} \geq 0$ for each $i$. Hence

$$
\langle v, v\rangle=\sum_{i=1}^{l} n_{i}\left\langle v, \alpha_{i}\right\rangle \leq 0 .
$$

It follows that $v=0$ and so $\Omega$ is empty.

Finally we must show that $\rho-\sum_{\alpha \in \Omega} \alpha$ either lies in a reflecting hyperplane or is a $W$-transform of $\rho$. Suppose it does not lie in any reflecting hyperplane. Then it lies in a chamber. Thus there exists $w \in W$ such that

$$
w\left(\rho-\sum_{\alpha \in \Omega} \alpha\right)
$$

lies in the fundamental chamber. However, $\rho-\sum_{\alpha \in \Omega} \alpha$ has the form $\frac{1}{2} \sum_{\alpha \in \Phi^{+}}( \pm \alpha)$, so $w\left(\rho-\sum_{\alpha \in \Omega} \alpha\right)$ also has form $\frac{1}{2} \sum_{\alpha \in \Phi^{+}}( \pm \alpha)$ since $w$ permutes the roots. Hence

$$
w\left(\rho-\sum_{\alpha \in \Omega} \alpha\right)=\rho-\sum_{\alpha \in \Omega^{\prime}} \alpha
$$

for some subset $\Omega^{\prime}$ of $\Phi^{+}$. Since this vector lies in the fundamental chamber, $\Omega^{\prime}$ must be empty. Hence

$$
w\left(\rho-\sum_{\alpha \in \Omega} \alpha\right)=\rho
$$

and so $\rho-\sum_{\alpha \in \Omega} \alpha$ is a $W$-transform of $\rho$
We can now prove Weyl's denominator formula.

Theorem 12.11 (Weyl's denominator formula).

$$
e_{\rho} \prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)=\sum_{w \in W} \varepsilon(w) e_{w(\rho)}
$$

where $\varepsilon(w)=(-1)^{l(w)}$.

Proof. Let $\mathbb{Z}\left[H^{*}\right]$ be the set of functions $f: H^{*} \rightarrow \mathbb{Z}$ of finite support. This is the set of finite $\mathbb{Z}$-combinations of the characteristic functions $e_{\lambda}$. Weyl's denominator formula is an identity in $\mathbb{Z}\left[H^{*}\right]$. There is a natural action of $W$ on $\mathbb{Z}\left[H^{*}\right]$ given by

$$
(w f) \lambda=f\left(w^{-1} \lambda\right)
$$

We define a map $\theta: \mathbb{Z}\left[H^{*}\right] \rightarrow \mathbb{Z}\left[H^{*}\right]$ by

$$
\theta(f)=\sum_{w \in W} \varepsilon(w) w f
$$

It is clear that, for $w^{\prime} \in W, \theta w^{\prime}=\varepsilon\left(w^{\prime}\right) \theta$, hence

$$
\theta\left(\varepsilon\left(w^{\prime}\right) w^{\prime}\right)=\theta
$$

and

$$
\theta^{2}=|W| \theta
$$

We now consider the effect of a fundamental reflection $s_{i}$ on

$$
\Delta=e_{\rho} \prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)
$$

We have

$$
s_{i}\left(e_{\rho} \prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)\right)=e_{s_{i}(\rho)} \prod_{\alpha \in \Phi^{+}}\left(1-e_{-s_{i}(\alpha)}\right)
$$

Now $s_{i}(\rho)=\rho-\alpha_{i}$ by the proof of Proposition 12.5. Also $s_{i}$ transforms every positive root to a positive root, except for $\alpha_{i}$. Hence we have

$$
\begin{aligned}
s_{i}\left(e_{\rho} \prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)\right) & =e_{\rho-\alpha_{i}}\left(\prod_{\substack{\alpha \in \Phi^{+} \\
\alpha \neq \alpha_{i}}}\left(1-e_{-\alpha}\right)\right)\left(1-e_{\alpha_{i}}\right) \\
& =e_{\rho}\left(\prod_{\substack{\alpha \in \Phi^{+} \\
\alpha \neq \alpha_{i}}}\left(1-e_{-\alpha}\right)\right)\left(e_{-\alpha_{i}}-1\right) \\
& =-e_{\rho} \prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)
\end{aligned}
$$

Thus $s_{i}(\Delta)=-\Delta$. It follows that $w(\Delta)=\varepsilon(w) \Delta$ for all $w \in W$. Hence $\theta(\Delta)=$ $|W| \Delta$.

We also have

$$
\begin{aligned}
\Delta & =e_{\rho} \prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right) \\
& =e_{\rho} \sum_{\Omega \subset \Phi^{+}}(-1)^{|\Omega|} e_{-\sum_{\alpha \in \Omega} \alpha} \\
& =\sum_{\Omega \subset \Phi^{+}}(-1)^{|\Omega|} e_{\rho-\sum_{\alpha \in \Omega} \alpha}
\end{aligned}
$$

Now $\rho-\sum_{\alpha \in \Omega} \alpha$ either is of form $w(\rho)$ for some $w \in W$ or lies in some reflecting hyperplane, by Proposition 12.10. If $v$ lies in a reflecting hyperplane then $\theta(v)=0$ since the terms in $\theta(v)$ cancel out in pairs. For if $v \in L_{\alpha}$ then

$$
\varepsilon\left(w w_{\alpha}\right) w w_{\alpha} v=-\varepsilon(w) w v
$$

Thus we have

$$
\theta(\Delta)=\theta \sum_{w \in W} \varepsilon(w) e_{w \rho}
$$

since if $\rho-\sum_{\alpha \in \Omega} \alpha=w(\rho)$ then $|\Omega|=l(w)$ by the proof of Proposition 12.10. Thus

$$
\theta(\Delta)=\theta\left(\theta\left(e_{\rho}\right)\right)=\theta^{2}\left(e_{\rho}\right)=|W| \theta\left(e_{\rho}\right) .
$$

But $\theta(\Delta)=|W| \Delta$ as shown above. It follows that

$$
\Delta=\theta\left(e_{\rho}\right)=\sum_{w \in W} \varepsilon(w) e_{w \rho}
$$

Corollary $12.12 \operatorname{ch} M(\lambda)=\frac{e_{\lambda+\rho}}{\sum_{w \in W} \varepsilon(w) e_{w \rho}}$.
Proof. This follows from Proposition 12.4 and Theorem 12.11.

### 12.4 Composition factors of Verma modules

We shall show in this section that each Verma module $M(\lambda)$ has a composition series of finite length and that all its composition factors are irreducible modules of the form $L(\mu)$ where $\mu=w(\lambda+\rho)-\rho$ for some $w \in W$. It will be convenient to define

$$
w \cdot \lambda=w(\lambda+\rho)-\rho
$$

We shall use these results in the following section to prove Weyl's character formula for the finite dimensional irreducible modules $L(\lambda)$.

We begin with a lemma on filtered algebras and their corresponding graded algebras. We recall the definitions as given in Section 11.1.

Lemma 12.13 Let $A=\bigcup_{i} A_{i}$ be a filtered algebra with

$$
A_{0} \subset A_{1} \subset A_{2} \subset \cdots
$$

and let $B=B_{0} \oplus B_{1} \oplus B_{2} \oplus \cdots$ be the corresponding graded algebra.
(i) if $I$ is a left ideal of $A$ then $\mathrm{gr} I=\bigoplus_{i} \frac{A_{i-1}+\left(A_{i} \cap I\right)}{A_{i-1}}$ is a left ideal of $B$.
(ii) If $I_{1} \subset I_{2}$ then $\operatorname{gr} I_{1} \subset \operatorname{gr} I_{2}$.
(iii) If $I_{1} \subset I_{2}$ and $\operatorname{gr} I_{1}=\operatorname{gr} I_{2}$ then $I_{1}=I_{2}$.
(iv) If $B$ satisfies the maximal condition on left ideals so does $A$.

Proof. We recall that $B=\bigoplus_{i} B_{i}$ where $B_{i}=A_{i} / A_{i-1}$. If $x \in A_{i} \cap I, y \in A_{j}$ then we have

$$
\left(A_{j-1}+y\right)\left(A_{i-1}+x\right)=A_{i+j-1}+y x
$$

where $y x \in A_{i+j} \cap I$. Thus $A_{i+j-1}+y x \in \operatorname{gr} I$. It follows that $\operatorname{gr} I$ is a left ideal of $B$.

It is clear from the definition that if $I_{1} \subset I_{2}$ then $\mathrm{gr} I_{1} \subset \mathrm{gr} I_{2}$.
We now suppose that $I_{1} \subset I_{2}$ and gr $I_{1}=\operatorname{gr} I_{2}$. Then

$$
A_{i-1}+\left(A_{i} \cap I_{1}\right)=A_{i-1}+\left(A_{i} \cap I_{2}\right)
$$

for each $i$. Thus we have

$$
A_{i} \cap I_{2}=\left(A_{i} \cap I_{1}\right)+\left(A_{i-1} \cap I_{2}\right)
$$

We shall show that $A_{i} \cap I_{1}=A_{i} \cap I_{2}$ by induction on $i$. We know $A_{0} \cap I_{1}=$ $A_{0} \cap I_{2}$ since $\operatorname{gr} I_{1}=\operatorname{gr} I_{2}$.

Assume inductively that $A_{i-1} \cap I_{1}=A_{i-1} \cap I_{2}$. Then we have

$$
A_{i} \cap I_{2}=\left(A_{i} \cap I_{1}\right)+\left(A_{i-1} \cap I_{1}\right)=A_{i} \cap I_{1}
$$

Thus $A_{i} \cap I_{2}=A_{i} \cap I_{1}$ for all $i$. Since $A=\cup_{i} A_{i}$ it follows that $I_{1}=I_{2}$.
Now suppose that

$$
I_{1} \subset I_{2} \subset I_{3} \subset \cdots
$$

is a chain of left ideals of $A$. Then

$$
\operatorname{gr} I_{1} \subset \operatorname{gr} I_{2} \subset \operatorname{gr} I_{3} \subset \cdots
$$

is a chain of left ideals of $B$. Assume that $B$ satisfies the maximal condition on left ideals. Then we have $\operatorname{gr} I_{i}=\operatorname{gr} I_{j}$ for all $i, j$ sufficiently large. It follows that $I_{i}=I_{j}$ for all $i, j$ sufficiently large. Hence $A$ satisfies the maximal condition on left ideals.

Proposition 12.14 $\mathfrak{U}(L)$ satisfies the maximal condition on left ideals.

Proof. $\mathfrak{U}(L)$ is a filtered algebra whose graded algebra is the symmetric algebra $S(L)$. However, $S(L)$ is isomorphic to the polynomial ring $\mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$ where $n=\operatorname{dim} L$, so satisfies the maximal condition on (left) ideals, by Hilbert's basis theorem. Thus $\mathfrak{H}(L)$ satisfies the maximal condition on left ideals, by Lemma 12.13.

Corollary 12.15 The Verma module $M(\lambda)$ satisfies the maximal condition on submodules.

Proof. The left ideals of $\mathfrak{H}(L)$ are the same as the $\mathfrak{H}(L)$-submodules. Thus $\mathfrak{U}(L)$ satisfies the maximal condition on submodules. We recall that

$$
M(\lambda)=\mathfrak{U}(L) / K_{\lambda}
$$

where $K_{\lambda}$ is a submodule of $\mathfrak{H}(L)$. It follows that $M(\lambda)$ satisfies the maximal condition on submodules.

Theorem 12.16 The Verma module $M(\lambda)$ has a finite composition series

$$
M(\lambda)=N_{0} \supset N_{1} \supset N_{2} \supset \cdots \supset N_{r}=O
$$

where each $N_{i}$ is a submodule of $M(\lambda)$ and $N_{i+1}$ is a maximal submodule of $N_{i}$. Moreover $N_{i} / N_{i+1}$ is isomorphic to $L(w \cdot \lambda)$ for some $w \in W$.

Proof. Since $M(\lambda)$ satisfies the maximal condition on submodules, every submodule of $M(\lambda)$ has a maximal submodule. Thus we have a descending series

$$
M(\lambda)=N_{0} \supset N_{1} \supset N_{2} \supset \cdots
$$

of submodules, in which $N_{i+1}$ is a maximal submodule of $N_{i}$. We wish to show that this series reaches $O$ after finitely many steps.

Now $M(\lambda)$ is the direct sum of its weight spaces by Theorem 10.7. Thus every submodule of $M(\lambda)$ is also the direct sum of its weight spaces, by the proof of Theorem 10.9. It follows that each quotient $N_{i} / N_{i+1}$ is the direct sum of its weight spaces. Moreover each weight $\mu$ of $N_{i} / N_{i+1}$ is a weight of $M(\lambda)$ so satisfies $\mu \prec \lambda$ with respect to the natural partial order on weights. Thus we can choose a weight $\mu$ of $N_{i} / N_{i+1}$ which is maximal in this partial order among the set of possible weights. Let $v$ be a non-zero vector in $N_{i} / N_{i+1}$ of weight $\mu$. Then we have $e_{i} v=0$ and $h v=\mu(h) v$ for all $h \in H$. Thus we have

$$
\mathfrak{U}(L) v=\mathfrak{l}\left(N^{-}\right) v
$$

However, $N_{i} / N_{i+1}$ is an irreducible $\mathfrak{U}(L)$-module, thus $\mathfrak{H}(L) v=N_{i} / N_{i+1}$. Thus we have a homomorphism $M(\mu) \rightarrow N_{i} / N_{i+1}$ given by $u m_{\mu} \rightarrow u v$ for all $u \in \mathfrak{H}\left(N^{-}\right)$as in Proposition 10.13. This homomorphism is surjective and its kernel is the unique maximal submodule of $M(\mu)$, since $N_{i} / N_{i+1}$ is irreducible. It follows that $N_{i} / N_{i+1}$ is isomorphic to $L(\mu)$, the unique irreducible quotient of $M(\mu)$.

We now consider the action of the centre $Z(L)$ of $\mathscr{H}(L) . Z(L)$ acts on $M(\lambda)$ by scalar multiplications. The element $z \in Z(L)$ acts on $M(\lambda)$ by scalar multiplication by $\chi_{\lambda}(z)$, as in Section 11.5. Hence $z$ acts on each submodule $N_{i}$ and each quotient $N_{i} / N_{i+1}$ as scalar multiplication by $\chi_{\lambda}(z)$. However,
$z$ acts on $M(\mu)$ as scalar multiplication by $\chi_{\mu}(z)$, and so also on its quotient $L(\mu)$. Since $N_{i} / N_{i+1}$ is isomorphic to $L(\mu)$ we deduce that $\chi_{\lambda}(z)=\chi_{\mu}(z)$ for all $z \in Z(L)$. Hence $\chi_{\lambda}=\chi_{\mu}$. It follows from Theorem 11.31 that $\mu+\rho=$ $w(\lambda+\rho)$ for some $w \in W$. This is equivalent to $\mu=w \cdot \lambda$ for some $w \in W$.

Now $W$ is finite and so there are only finitely many possible composition factors of $M(\lambda)$, up to isomorphism. Also each weight space of $M(\lambda)$ is finite dimensional. Thus $L(\mu)$, which contains $\mu$ as a weight, can appear as a composition factor with multiplicity at most the dimension of the $\mu$ weight space $M(\lambda)_{\mu}$. It follows that the series

$$
M(\lambda)=N_{0} \supset N_{1} \supset N_{2} \supset \cdots
$$

must reach $O$ after at most $\sum_{w \in W} \operatorname{dim} M(\lambda)_{w \cdot \lambda}$ steps. Thus $M(\lambda)$ has a finite composition series and each composition factor has form $L(w \cdot \lambda)$ for some $w \in W$.

### 12.5 Weyl's character formula

We now find a formula for the characters of the finite dimensional irreducible modules $L(\lambda)$ where $\lambda \in X^{+}$.

Theorem 12.17 (Weyl's character formula). Let $\lambda \in X^{+}$. Then

$$
\operatorname{ch} L(\lambda)=\frac{\sum_{w \in W} \varepsilon(w) e_{w(\lambda+\rho)}}{\sum_{w \in W} \varepsilon(w) e_{w(\rho)}}
$$

(This is an equality in the ring $\mathfrak{R}$ of Section 12.1 since the denominator

$$
\Delta=\sum_{w \in W} \varepsilon(w) e_{w(\rho)}
$$

is an invertible element of $\mathfrak{R}$.)

Proof. Since $\lambda$ is a dominant integral weight we have $\lambda\left(h_{i}\right) \geq 0$ for $i=1, \ldots, l$. Hence $(\lambda+\rho)\left(h_{i}\right)=\lambda\left(h_{i}\right)+1>0$ for $i=1, \ldots, l$. Thus $\lambda+\rho$ lies in the fundamental chamber $C$. Hence $w(\lambda+\rho)$ lies in the chamber $w(C)$. It follows from Proposition 12.7 that the weights $w(\lambda+\rho)$ for $w \in W$ are all distinct.

Now the highest weight of the Verma module $M(w \cdot \lambda)$ is $w \cdot \lambda=w(\lambda+\rho)-\rho$. Thus the characters ch $M(w \cdot \lambda) \in \mathfrak{R}$ are linearly independent as $w$ runs over $W$. Similarly the characters ch $L(w \cdot \lambda)$ are linearly independent for $w \in W$.

Now $M(w \cdot \lambda)$ has a finite composition series with composition factors of form $L(y \cdot \lambda)$ for $y \in W$, by Theorem 12.16. Moreover, since $y \cdot \lambda$ is a weight of $L(y \cdot \lambda)$ and $w \cdot \lambda$ is the highest weight of $M(w \cdot \lambda)$ we have $y \cdot \lambda \prec w \cdot \lambda$ whenever $L(y \cdot \lambda)$ occurs as a composition factor of $M(w \cdot \lambda)$. Moreover $w \cdot \lambda$ occurs as a weight of $M(w \cdot \lambda)$ with multiplicity 1 , thus $L(w \cdot \lambda)$ appears as a composition factor of $M(w \cdot \lambda)$ with multiplicity 1 . We therefore have

$$
\operatorname{ch} M(w \cdot \lambda)=\sum_{y \in W} a_{w y} \operatorname{ch} L(y \cdot \lambda)
$$

where $a_{w y} \in \mathbb{Z}, a_{w y} \geq 0, a_{w w}=1$, and $a_{w y} \neq 0$ only if $y \cdot \lambda \prec w \cdot \lambda$. If we write the elements of $W$ in an order compatible with the partial order $y \cdot \lambda \prec w \cdot \lambda$ we see that the integers $a_{w y}$ form a triangular $|W| \times|W|$ matrix with entries 1 on the diagonal. The determinant of this matrix is 1 . Thus we may invert the above equations to obtain

$$
\operatorname{ch} L(w \cdot \lambda)=\sum_{y \in W} b_{w y} \operatorname{ch} M(y \cdot \lambda)
$$

where $b_{w y} \in \mathbb{Z}$ and $b_{w w}=1$. (The $b_{w y}$ will no longer be non-negative.) In particular we have

$$
\operatorname{ch} L(\lambda)=\sum_{y \in W} c_{y} \operatorname{ch} M(y \cdot \lambda)
$$

where $c_{y}=b_{1 y}$. By Proposition 12.4 this gives

$$
\operatorname{ch} L(\lambda)=\frac{\sum_{y \in W} c_{y} e_{y(\lambda+\rho)}}{\Delta}
$$

where $c_{1}=1$. We wish to determine the remaining coefficients $c_{y}$.
We recall from Proposition 10.22 that

$$
\operatorname{dim} L(\lambda)_{\mu}=\operatorname{dim} L(\lambda)_{w(\mu)}
$$

for all $w \in W$. Thus we have

$$
w(\operatorname{ch} L(\lambda))=\operatorname{ch} L(\lambda) \quad \text { for all } w \in W
$$

On the other hand we have

$$
s_{i}(\Delta)=-\Delta
$$

by Theorem 12.11 and thus

$$
w(\Delta)=\varepsilon(w) \Delta
$$

It follows that

$$
w\left(\sum_{y \in W} c_{y} e_{y(\lambda+\rho)}\right)=\varepsilon(w)\left(\sum_{y \in W} c_{y} e_{y(\lambda+\rho)}\right) .
$$

Thus

$$
\sum_{y \in W} c_{y} e_{w y(\lambda+\rho)}=\sum_{y \in W} \varepsilon(w) c_{y} e_{y(\lambda+\rho)}
$$

since $w e_{\lambda}=e_{w \lambda}$. This is equivalent to

$$
\sum_{y \in W} c_{w^{-1} y} e_{y(\lambda+\rho)}=\sum_{y \in W} \varepsilon(w) c_{y} e_{y(\lambda+\rho)} .
$$

Since the functions $e_{y(\lambda+\rho)}$ for $y \in W$ are linearly independent we deduce that

$$
c_{w^{-1} y}=\varepsilon(w) c_{y} .
$$

In particular we have $c_{w^{-1}}=\varepsilon(w)$, thus $c_{w}=\varepsilon\left(w^{-1}\right)=\varepsilon(w)$. It follows that

$$
\operatorname{ch} L(\lambda)=\frac{\sum_{w \in W} \varepsilon(w) e_{w(\lambda+\rho)}}{\Delta}
$$

as required.
We note that in the special case $\lambda=0$ we have $\operatorname{ch} L(0)=e_{0}$ as $L(0)$ is the trivial 1-dimensional representation of $L$. Thus we have

$$
\sum_{w \in W} \varepsilon(w) e_{w(\rho)}=\Delta e_{0}=\Delta .
$$

This gives an alternative proof of Weyl's denominator formula, Theorem 12.11.

We also note that while the character ch $L(\lambda)$ is invariant under the Weyl group both the numerator and the denominator in Weyl's character formula are alternating functions under the Weyl group, i.e. satisfy $w(a)=\varepsilon(w) a$.

We may deduce from Weyl's character formula a formula due to Kostant for the dimension of the weight space $L(\lambda)_{\mu}$ of $L(\lambda)$.

Theorem 12.18 (Kostant's multiplicity formula). Let $\lambda \in X^{+}$and $\mu \in X$. Then

$$
\operatorname{dim} L(\lambda)_{\mu}=\sum_{w \in W} \varepsilon(w) \mathfrak{W}(w(\lambda+\rho)-(\mu+\rho))
$$

where $\mathfrak{B}$ is the partition function defined in Theorem 10.7.
Proof. We have ch $L(\lambda)=\sum_{\mu} \operatorname{dim} L(\lambda)_{\mu} e_{\mu}$. Moreover we know that

$$
\Delta^{-1}=e_{-\rho} \Gamma=e_{-\rho} \sum_{\nu} \mathfrak{\beta}(\nu) e_{-\nu}
$$

by Lemma 12.3. Thus Weyl's character formula gives the identity

$$
\begin{aligned}
\sum_{\mu} \operatorname{dim} L(\lambda)_{\mu} e_{\mu} & =\left(\sum_{w \in W} \varepsilon(w) e_{w(\lambda+\rho)}\right) e_{-\rho} \sum_{\nu} \mathfrak{P}(\nu) e_{-\nu} \\
& =\sum_{w \in W} \sum_{\nu} \varepsilon(w) \mathfrak{P}(\nu) e_{w(\lambda+\rho)-\rho-\nu}
\end{aligned}
$$

We compare the coefficients of $e_{\mu}$ on both sides. This gives

$$
\operatorname{dim} L(\lambda)_{\mu}=\sum_{w \in W} \varepsilon(w) \mathfrak{P}(w(\lambda+\rho)-(\mu+\rho))
$$

We can also derive from Weyl's character formula a formula for the dimension of $L(\lambda)$.

Theorem 12.19 (Weyl's dimension formula). Let $\lambda \in X^{+}$. Then

$$
\operatorname{dim} L(\lambda)=\frac{\prod_{\alpha \in \Phi^{+}}\langle\lambda+\rho, \alpha\rangle}{\prod_{\alpha \in \Phi^{+}}\langle\rho, \alpha\rangle}
$$

Proof. Let $\mathfrak{R}_{0}$ be the subring of $\mathfrak{R}$ consisting of all finite sums $\sum_{\mu \in X} n_{\mu} e_{\mu}$ with $n_{\mu} \in \mathbb{Z}$. Then the character formula

$$
\Delta \operatorname{ch} L(\lambda)=\sum_{w \in W} \varepsilon(w) e_{w(\lambda+\rho)}
$$

may be regarded as an identity in $\mathfrak{R}_{0}$. Let $A=\mathbb{R}[[t]]$ be the ring of formal power series in the variable $t$ with real coefficients. Then for each weight $\xi \in X$ we have a ring homomorphism

$$
\theta_{\xi}: \Re_{0} \rightarrow A
$$

given by

$$
\theta_{\xi}\left(e_{\mu}\right)=\exp (\langle\xi, \mu\rangle t)=1+\langle\xi, \mu\rangle t+\frac{1}{2!}\langle\xi, \mu\rangle^{2} t^{2}+\cdots
$$

Consider $\theta_{\xi}\left(\sum_{w \in W} \varepsilon(w) e_{w \mu}\right)$. We have

$$
\begin{aligned}
\theta_{\xi}\left(\sum_{w \in W} \varepsilon(w) e_{w \mu}\right) & =\sum_{w \in W} \varepsilon(w) \exp (\langle\xi, w \mu\rangle t)=\sum_{w \in W} \varepsilon(w) \exp \left(\left\langle\mu, w^{-1} \xi\right\rangle t\right) \\
& =\sum_{w \in W} \varepsilon(w) \exp (\langle\mu, w \xi\rangle t)=\theta_{\mu}\left(\sum_{w \in W} \varepsilon(w) e_{w \xi}\right)
\end{aligned}
$$

In particular we have

$$
\begin{aligned}
\theta_{\rho}\left(\sum_{w \in W} \varepsilon(w) e_{w(\lambda+\rho)}\right) & =\theta_{\lambda+\rho}\left(\sum_{w \in W} \varepsilon(w) e_{w \rho}\right) \\
& =\theta_{\lambda+\rho}\left(e_{-\rho} \prod_{\alpha \in \Phi^{+}}\left(e_{\alpha}-1\right)\right) \\
& =\exp (\langle\lambda+\rho,-\rho\rangle t) \prod_{\alpha \in \Phi^{+}}(\exp \langle\lambda+\rho, \alpha\rangle t-1) \\
& =\exp (\langle\lambda+\rho,-\rho\rangle t) \prod_{\alpha \in \Phi^{+}}(\langle\lambda+\rho, \alpha\rangle t+\cdots) \\
& =t^{N}\left(\prod_{\alpha \in \Phi^{+}}\langle\lambda+\rho, \alpha\rangle+\cdots\right) \quad \text { where } N=\left|\Phi^{+}\right|
\end{aligned}
$$

By putting $\lambda=0$ we obtain

$$
\theta_{\rho}\left(\sum_{w \in W} \varepsilon(w) e_{w \rho}\right)=t^{N}\left(\prod_{\alpha \in \Phi^{+}}\langle\rho, \alpha\rangle+\cdots\right) .
$$

Thus by applying $\theta_{\rho}$ to Weyl's character formula we obtain

$$
t^{N}\left(\prod_{\alpha \in \Phi^{+}}\langle\rho, \alpha\rangle+\cdots\right) \sum_{\mu} \operatorname{dim} L(\lambda)_{\mu} \exp (\langle\rho, \mu\rangle t)=t^{N}\left(\prod_{\alpha \in \Phi^{+}}\langle\lambda+\rho, \alpha\rangle+\cdots\right) .
$$

By cancelling $t^{N}$ and then taking the constant term we obtain

$$
\left(\prod_{\alpha \in \Phi^{+}}\langle\rho, \alpha\rangle\right) \operatorname{dim} L(\lambda)=\prod_{\alpha \in \Phi^{+}}\langle\lambda+\rho, \alpha\rangle
$$

### 12.6 Complete reducibility

We have now attained a good understanding of the finite dimensional irreducible modules for a semisimple Lie algebra $L$. We now consider arbitrary finite dimensional $L$-modules. Each of these turns out to be a direct sum of irreducible $L$-modules.

Theorem 12.20 Let $L$ be a semisimple Lie algebra and $V$ a finite dimensional $L$-module. Then $V$ is completely reducible.

Proof. We shall prove this result in a number of steps. If $V$ is itself irreducible there is nothing to prove. Thus we suppose $U$ is a proper submodule of $V$. It
will be sufficient to show that $U$ has a complementary submodule $U^{\prime}$ in $V$, that is a submodule such that $V=U \oplus U^{\prime}$.
(a) Suppose $\operatorname{dim} V=2, \operatorname{dim} U=1$. Then $U$ and $V / U$ are 1-dimensional $L$-modules. Since for $x, y \in L, u \in U$ we have

$$
[x y] u=x(y u)-y(x u)
$$

and since the actions of $x$ and $y$ on the 1-dimensional module commute we have

$$
[x y] u=0
$$

Thus [LL] acts as 0 on $U$. Since $L$ is semisimple we have $[L L]=L$. Thus $U$ gives the trivial 1-dimensional representation $L(0)$. Similarly $V / U$ is isomorphic to $L(0)$.

Now let $v \in V$. Then

$$
[x y] v=x(y v)-y(x v)
$$

Since $L$ annihilates $V / U$ we have $x v \in U$ and $y v \in U$. Since $L$ annihilates $U$ we have $x(y v)=0$ and $y(x v)=0$. Hence $[x y] v=0$. This shows that [ $L L$ ] annihilates $V$, i.e. $L$ annihilates $V$. But then any complementary subspace $U^{\prime}$ of $U$ is a submodule of $V$.
(b) Suppose $U$ is irreducible, $\operatorname{dim} U>1$, and $\operatorname{dim}(V / U)=1$.

Then $U$ is isomorphic to $L(\lambda)$ for some $\lambda \in X^{+}$with $\lambda \neq 0$. We consider the action of the Casimir element $c$ on $V$. We recall from Proposition 11.36 that $c$ acts on the irreducible module $L(\lambda)$ as scalar multiplication by $\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle$. In particular $c$ acts on $L(0)$ as zero, and $c$ acts on $L(\lambda)$ for $\lambda \in X^{+}, \lambda \neq 0$, as multiplication by a positive scalar. For then $\langle\lambda, \lambda\rangle>0$ and $\langle\lambda, \rho\rangle \geq 0$ since $\lambda \in X^{+}$. Thus $c$ has one eigenvalue 0 on $V$ and $\operatorname{dim} V-1$ eigenvalues $\chi_{\lambda}(c)=\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle>0$. Let $U^{\prime}$ be the eigenspace of $c$ on $V$ with eigenvalue 0 . Then we have

$$
V=U \oplus U^{\prime}
$$

Moreover $U^{\prime}$ is a submodule of $V$. For let $x \in L, u^{\prime} \in U^{\prime}$. Then

$$
c\left(x u^{\prime}\right)=x\left(c u^{\prime}\right)=0
$$

since $c$ lies in the centre $Z(L)$ of $\mathfrak{H}(L)$. Thus $x u^{\prime} \in U^{\prime}$ and $U^{\prime}$ is the required submodule of $V$.
(c) Suppose $\operatorname{dim}(V / U)=1$ but $U$ is not irreducible. We prove the existence of the required complementary submodule $U^{\prime}$ by induction on $\operatorname{dim} U$. Let $\bar{U}$ be a proper submodule of $U$. Then by induction we have

$$
V / \bar{U}=U / \bar{U} \oplus \bar{V} / \bar{U}
$$

for some submodule $\bar{V}$ of $V$ containing $\bar{U}$. We have $\operatorname{dim}(\bar{V} / \bar{U})=1$ and $\operatorname{dim} \bar{U}<\operatorname{dim} U$. Thus we may apply induction again and conclude that there exists a submodule $U^{\prime}$ such that

$$
\bar{V}=\bar{U} \oplus U^{\prime}
$$

But then we have $V=U \oplus U^{\prime}$ as required
(d) We now consider the general case when $U$ is any proper submodule of $V$. We consider the set $\operatorname{Hom}(V, U)$ of all linear maps from $V$ to $U$. We can make this set into an $L$-module as follows. If $x \in L$ and $\theta \in \operatorname{Hom}(V, U)$ we define $x \theta \in \operatorname{Hom}(V, U)$ by

$$
(x \theta) v=x(\theta(v))-\theta(x v) \in U
$$

Then we have

$$
\begin{aligned}
(y(x \theta)) v & =y((x \theta) v)-(x \theta)(y v) \\
& =y(x(\theta v))-y(\theta(x v))-x(\theta(y v))+\theta(x(y v))
\end{aligned}
$$

Similarly

$$
(x(y \theta)) v=x(y(\theta v))-x(\theta(y v))-y(\theta(x v))+\theta(y(x v)) .
$$

Thus

$$
\begin{aligned}
(x(y \theta)-y(x \theta)) v & =x(y(\theta v))-y(x(\theta v))+\theta(y(x v))-\theta(x(y v)) \\
& =[x y](\theta v)-\theta([x y] v) \\
& =([x y] \theta) v .
\end{aligned}
$$

Thus $\operatorname{Hom}(V, U)$ is an $L$-module. Let $S$ be the subspace of $\operatorname{Hom}(V, U)$ of maps $\theta$ such that $\left.\theta\right|_{U}$ is a scalar multiplication. Then $S$ is a submodule of $\operatorname{Hom}(V, U)$. For suppose $\theta \in S, x \in L$. Then for $u \in U$ we have

$$
(x \theta) u=x(\theta u)-\theta(x u)=\xi x u-\xi x u=0
$$

where $\theta$ acts on $U$ as multiplication by $\xi$. Thus $S$ is a submodule of $\operatorname{Hom}(V, U)$. Moreover let $T$ be the subspace of $S$ of maps $\theta$ such that $\left.\theta\right|_{U}$ is zero. Then $T$ is a submodule of $S$ and $\operatorname{dim}(S / T)=1$.

We know then from the earlier parts of the proof that there is a submodule $T^{\prime}$ of $S$ such that $S=T \oplus T^{\prime}$. We have $\operatorname{dim} T^{\prime}=1$. Suppose $T^{\prime}$ is
spanned by the non-zero element $f: V \rightarrow U$. We may choose $f$ so that $f(u)=u$ for all $u \in U$. We have $x f=0$ for all $x \in L$ since $\operatorname{dim} T^{\prime}=1$. Thus

$$
(x f) v=x(f v)-f(x v)=0
$$

for all $v \in V$, that is

$$
x(f v)=f(x v) \quad \text { for all } x \in L, v \in V
$$

This shows that $f: V \rightarrow U$ is a homomorphism of $L$-modules. Let $U^{\prime}$ be the kernel of $f$. Then $U^{\prime}$ is a submodule of $V$. We have $U \cap U^{\prime}=O$ since $f$ acts as the identity on $U$, and

$$
\operatorname{dim} V=\operatorname{dim} U+\operatorname{dim} U^{\prime}
$$

since $f$ is surjective. Hence we have

$$
V=U \oplus U^{\prime}
$$

and $U^{\prime}$ is the required complementary submodule.
Note The crucial step in the above proof of complete reducibility is the fact that the Casimir element $c$ acts on the irreducible module $L(\lambda)$ for $\lambda \in X^{+}, \lambda \neq 0$, as multiplication by a positive scalar.

The theorem of complete reducibility shows that every finite dimensional $L$-module is a direct sum of irreducible $L$-modules each isomorphic to $L(\lambda)$ for some $\lambda \in X^{+}$.

In particular the tensor product $L(\lambda) \otimes L(\mu)$ is a direct sum of irreducible modules $L(\nu)$ where $\lambda, \mu, v \in X^{+}$. It is natural to try to determine the multiplicity with which $L(\nu)$ occurs as a direct summand of $L(\lambda) \otimes L(\mu)$. This multiplicity is given in a formula of Steinberg.

Theorem 12.21 (Steinberg's multiplicity formula). Let $\lambda, \mu \in X^{+}$and

$$
L(\lambda) \otimes L(\mu)=\sum_{\nu \in X^{+}} c_{\lambda \mu \nu} L(\nu)
$$

Then

$$
c_{\lambda \mu \nu}=\sum_{w, w^{\prime} \in W} \varepsilon(w) \varepsilon\left(w^{\prime}\right) \mathfrak{B}\left(w(\lambda+\rho)+w^{\prime}(\mu+\rho)-(\nu+2 \rho)\right)
$$

Proof. We have

$$
\operatorname{ch} L(\lambda) \operatorname{ch} L(\mu)=\sum_{\nu \in X^{+}} c_{\lambda \mu \nu} \operatorname{ch} L(\nu)
$$

by Lemma 12.2. We multiply both sides of this equation by the Weyl denominator $\Delta$. By Weyl's character formula, Theorem 12.17 , we have

$$
\left(\sum_{w \in W} \varepsilon(w) e_{w(\lambda+\rho)}\right)\left(\sum_{\xi \in X} \operatorname{dim} L(\mu)_{\xi} e_{\xi}\right)=\sum_{\nu \in X^{+}} c_{\lambda \mu \nu}\left(\sum_{w \in W} \varepsilon(w) e_{w(\nu+\rho)}\right)
$$

Thus $\quad \sum_{w \in W} \sum_{\xi \in X} \varepsilon(w) \operatorname{dim} L(\mu)_{\xi} e_{w(\lambda+\rho)+\xi}=\sum_{\nu \in X^{+}} \sum_{w \in W} c_{\lambda \mu \nu} \varepsilon(w) e_{w(\nu+\rho)}$. Now $\nu \in X^{+}$, thus $\nu+\rho$ lies in the fundamental chamber $C$. Thus $w(\nu+\rho)$ lies in the chamber $w(C)$. Thus the elements $w(\nu+\rho)$ are all distinct as $w, \nu$ vary, and so the elements $e_{w(\nu+\rho)}$ are linearly independent. We may therefore compare the coefficients of $e_{w(\nu+\rho)}$ on both sides of the above equation. In fact we compare the coefficients of $e_{\nu+\rho}$ on both sides. This gives

$$
\begin{aligned}
c_{\lambda \mu \nu} & =\sum_{\substack{w \in W \\
w(\lambda+\rho)+\xi=\nu+\rho}} \sum_{\xi \in X} \varepsilon(w) \operatorname{dim} L(\mu)_{\xi} \\
& =\sum_{w \in W} \varepsilon(w) \operatorname{dim} L(\mu)_{\nu+\rho-w(\lambda+\rho)}
\end{aligned}
$$

We now use Kostant's multiplicity formula, Theorem 12.18. This gives

$$
\operatorname{dim} L(\mu)_{\nu+\rho-w(\lambda+\rho)}=\sum_{w^{\prime} \in W} \varepsilon\left(w^{\prime}\right) \mathfrak{P}\left(w^{\prime}(\mu+\rho)+w(\lambda+\rho)-(\nu+2 \rho)\right) .
$$

Thus we obtain

$$
c_{\lambda \mu \nu}=\sum_{w, w^{\prime} \in W} \varepsilon(w) \varepsilon\left(w^{\prime}\right) \mathfrak{B}\left(w(\lambda+\rho)+w^{\prime}(\mu+\rho)-(\nu+2 \rho)\right) .
$$

## 13

## Fundamental modules for simple Lie algebras

### 13.1 An alternative form of Weyl's dimension formula

Let $L$ be a finite dimensional simple Lie algebra. The irreducible $L$-modules $L\left(\omega_{i}\right)$ whose highest weights are the fundamental weights $\omega_{1}, \ldots, \omega_{l}$ are called the fundamental modules. In this chapter we shall determine the dimensions of the fundamental modules for the various simple Lie algebras. We shall first derive an alternative form of the Weyl dimension formula which will be useful in this respect.

Theorem 13.1 Let $\lambda=\sum_{i=1}^{l} m_{i} \omega_{i}$ be a dominant integral weight. Then

$$
\operatorname{dim} L(\lambda)=\prod_{\alpha \in \Phi^{+}} d_{\alpha}
$$

where $\alpha=\sum_{i=1}^{l} k_{i} \alpha_{i}$ and

$$
d_{\alpha}=\frac{\sum_{i}\left(m_{i}+1\right) k_{i} w_{i}}{\sum_{i} k_{i} w_{i}} .
$$

Here the integer $w_{i}$, called the weight of $\alpha_{i}$, is defined by

$$
\left\langle\alpha_{i}, \alpha_{i}\right\rangle=w_{i}\left\langle\alpha_{i_{0}}, \alpha_{i_{0}}\right\rangle
$$

where $\alpha_{i_{0}}$ is a short fundamental root. Thus $w_{i} \in\{1,2,3\}$ for each $i$.
Proof. We know from Theorem 12.19 that

$$
\operatorname{dim} L(\lambda)=\prod_{\alpha \in \Phi^{+}} d_{\alpha}
$$

where $d_{\alpha}=\frac{\langle\lambda+\rho, \alpha\rangle}{\langle\rho, \alpha\rangle}$.

Since $\lambda=\sum m_{i} \omega_{i}, \rho=\sum \omega_{i}, \alpha=\sum k_{i} \alpha_{i}$ we have

$$
d_{\alpha}=\frac{\left\langle\sum\left(m_{i}+1\right) \omega_{i}, \sum k_{i} \alpha_{i}\right\rangle}{\left\langle\sum \omega_{i}, \sum k_{i} \alpha_{i}\right\rangle}
$$

Now we know from the proof of Proposition 10.18 that $\left\langle\omega_{i}, \alpha_{j}\right\rangle=0$ if $i \neq j$ and $\left\langle\omega_{i}, \alpha_{i}\right\rangle=\frac{1}{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle$. Thus

$$
\begin{aligned}
d_{\alpha} & =\frac{\sum_{i}\left(m_{i}+1\right) k_{i} \cdot \frac{1}{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\sum_{i} k_{i} \cdot \frac{1}{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle} \\
& =\frac{\sum_{i}\left(m_{i}+1\right) k_{i} w_{i}}{\sum_{i} k_{i} w_{i}}
\end{aligned}
$$

### 13.2 Fundamental modules for $A_{l}$

The fundamental weights $\omega_{1}, \ldots, \omega_{l}$ for a simple Lie algebra of type $A_{l}$ will be numbered according to the vertices of the Dynkin diagram as shown.


We shall use Theorem 13.1 to calculate $\operatorname{dim} L\left(\omega_{j}\right)$ for $j \in\{1, \ldots, l\}$. We have $\operatorname{dim} L\left(\omega_{j}\right)=\prod_{\alpha \in \Phi^{+}} d_{\alpha}$ where

$$
d_{\alpha}=\frac{\sum_{i}\left(m_{i}+1\right) k_{i}}{\sum_{i} k_{i}} .
$$

(All weights $w_{i}$ are equal to 1 .) Now $m_{j}=1$ and $m_{i}=0$ if $i \neq j$. Thus if $\alpha$ does not involve the fundamental root $\alpha_{j}$ we have $d_{\alpha}=1$.

So suppose $\alpha$ does involve $\alpha_{j}$. Then $\alpha=\alpha_{i}+\cdots+\alpha_{j}+\cdots+\alpha_{k}$ for some $i$ with $1 \leq i \leq j$ and some $k$ with $j \leq k \leq l$. For such a root $\alpha$ we have

$$
d_{\alpha}=\frac{k-i+2}{k-i+1}
$$

Thus

$$
\begin{aligned}
\operatorname{dim} L\left(\omega_{j}\right) & =\prod_{\substack{i \\
1 \leq i \leq j}} \prod_{\substack{k \\
j \leq k \leq l}} \frac{k-i+2}{k-i+1} \\
& =\prod_{\substack{k \\
j \leq k \leq l}} \frac{(k+1) k \ldots(k-j+2)}{k(k-1) \ldots(k-j+1)}=\prod_{\substack{k \\
j \leq k \leq l}} \frac{k+1}{k-j+1} \\
& =\frac{(j+1)(j+2) \ldots(l+1)}{1.2 \ldots \ldots(l+1-j)}=\frac{(l+1)!}{j!(l+1-j)!}=\binom{l+1}{j}
\end{aligned}
$$

Thus we have shown

Proposition 13.2 The dimensions of the fundamental modules for the simple Lie algebra of type $A_{l}$ are

$l+1 \quad\binom{l+1}{2}\binom{l+1}{3} \quad\binom{l+1}{3}\binom{l+1}{2} \quad l+1$

These fundamental modules may be described in terms of exterior powers of the natural $A_{l}$-module of dimension $l+1$. We recall from Theorem 8.1 that $A_{l}$ is isomorphic to the Lie algebra $\mathfrak{S l}_{l+1}(\mathbb{C})$ of all $(l+1) \times(l+1)$ matrices of trace 0 . The identity map gives an $(l+1)$-dimensional representation of $A_{l}$ called the natural representation. Its weights are the maps $\mu_{1}, \mu_{2}, \ldots, \mu_{l+1}$ given by

$$
\left(\begin{array}{ccccc}
\lambda_{1} & & & & \\
& \lambda_{2} & & & \\
& & \cdot & & \\
& & & \cdot & \\
& & & & \\
& & & & \lambda_{l+1}
\end{array}\right) \xrightarrow{\mu_{i}} \lambda_{i}
$$

Then we have

$$
\begin{aligned}
& \mu_{1}-\mu_{2}=\alpha_{1} \\
& \vdots \\
& \mu_{l}-\mu_{l+1}=\alpha_{l} \\
& \mu_{1}+\cdots+\mu_{l+1}=0 .
\end{aligned}
$$

On the other hand we have $\alpha_{i}=\sum_{j} A_{j i} \omega_{j}$ by Proposition 10.17. Hence

$$
\begin{aligned}
& \alpha_{1}=2 \omega_{1}-\omega_{2} \\
& \alpha_{2}=-\omega_{1}+2 \omega_{2}-\omega_{3} \\
& \vdots \\
& \alpha_{l-1}=-\omega_{l-2}+2 \omega_{l-1}-\omega_{l} \\
& \alpha_{l}=\quad-\omega_{l-1}+2 \omega_{l} .
\end{aligned}
$$

Eliminating $\alpha_{1}, \ldots, \alpha_{l}$ we obtain

$$
\begin{aligned}
& \mu_{1}=\omega_{1} \\
& \mu_{2}=-\omega_{1}+\omega_{2} \\
& \vdots \\
& \mu_{l}=-\omega_{l-1}+\omega_{l} \\
& \mu_{l+1}=-\omega_{l} .
\end{aligned}
$$

Now the weights of the natural module satisfy $\mu_{1} \succ \mu_{2} \succ \cdots \succ \mu_{l+1}$ since $\mu_{i}-\mu_{i+1}=\alpha_{i}$. Thus the highest weight of the natural module is $\mu_{1}=\omega_{1}$. It follows that $L\left(\omega_{1}\right)$ is one of the irreducible direct summands of the natural module $V$. Since

$$
\operatorname{dim} V=\operatorname{dim} L\left(\omega_{1}\right)=l+1
$$

it follows that $V=L\left(\omega_{1}\right)$. Thus we have shown
Proposition 13.3 The natural $A_{l}$-module is an irreducible module with highest weight $\omega_{1}$.

To obtain the remaining fundamental $A_{l}$-modules we introduce exterior powers of modules.

### 13.3 Exterior powers of modules

Let $V$ be a finite dimensional module for a Lie algebra $L$. Let

$$
T(V)=T^{0}(V) \oplus T^{1}(V) \oplus T^{2}(V) \oplus \cdots
$$

be the tensor algebra of $V$, where

$$
T^{n}(V)=V \otimes \cdots \otimes V \quad(n \text { factors })
$$

$T(V)$ may be made into an associative algebra in which

$$
\left(x_{1} \otimes \cdots \otimes x_{m}\right)\left(y_{1} \otimes \cdots \otimes y_{n}\right)=x_{1} \otimes \cdots \otimes x_{m} \otimes y_{1} \otimes \cdots \otimes y_{n}
$$

for $x_{1}, \ldots, x_{m}, y_{1}, \ldots, y_{n} \in V$.
$T(V)$ may also be given the structure of an $L$-module satisfying

$$
x\left(x_{1} \otimes \cdots \otimes x_{m}\right)=\sum_{i=1}^{m} x_{1} \otimes \cdots \otimes x_{i-1} \otimes x x_{i} \otimes x_{i+1} \otimes \cdots \otimes x_{m}
$$

for all $x \in L$.

Let $J$ be the 2-sided ideal of $T(V)$ generated by the elements $v \otimes v$ for all $v \in V$.

Definition 13.4 $\Lambda(V)=T(V) / J$ is called the exterior algebra of $V$.
Let $v, v^{\prime} \in V$. Then

$$
\left(v+v^{\prime}\right) \otimes\left(v+v^{\prime}\right)=v \otimes v+v^{\prime} \otimes v^{\prime}+v \otimes v^{\prime}+v^{\prime} \otimes v
$$

Hence

$$
v \otimes v^{\prime}+v^{\prime} \otimes v \in J
$$

Now let $v_{1}, \ldots, v_{n}$ be a basis of $V$. Then $J$ is the 2 -sided ideal of $T(V)$ generated by all elements of form

$$
\begin{aligned}
& v_{i} \otimes v_{i} \quad i=1, \ldots, n \\
& v_{i} \otimes v_{j}+v_{j} \otimes v_{i} \quad i<j .
\end{aligned}
$$

It follows from this that

$$
J=\bigoplus_{k \geq 0}\left(T^{k}(V) \cap J\right)
$$

and that $T^{0}(V) \cap J=O, T^{1}(V) \cap J=O$. Hence

$$
\Lambda(V)=\Lambda^{0} V \oplus \Lambda^{1} V \oplus \Lambda^{2} V \oplus \cdots
$$

where $\Lambda^{k} V=T^{k}(V) / T^{k}(V) \cap J$. In particular we have

$$
\begin{aligned}
& \Lambda^{0} V \cong T^{0}(V)=\mathbb{C} 1 \\
& \Lambda^{1} V \cong T^{1}(V)=V
\end{aligned}
$$

Thus we may identify the subspace $\Lambda^{0} V \oplus \Lambda^{1} V$ of $\Lambda(V)$ with $\mathbb{C} 1 \oplus V$.
Let $\sigma: T(V) \rightarrow \Lambda(V)=T(V) / J$ be the natural homomorphism. We define

$$
\sigma\left(v \otimes v^{\prime}\right)=v \wedge v^{\prime} \quad \text { for } v, v^{\prime} \in V
$$

Then every element of $\Lambda(V)$ is a linear combination of elements

$$
v_{i_{1}} \wedge \cdots \wedge v_{i_{k}} \quad i_{1}, \ldots, i_{k} \in\{1, \ldots, n\} .
$$

The relations defining $J$ may be written

$$
\begin{aligned}
& v_{i} \wedge v_{i}=0 \quad i \in\{1, \ldots, n\} \\
& v_{j} \wedge v_{i}=-\left(v_{i} \wedge v_{j}\right) \quad i<j
\end{aligned}
$$

By applying these relations we see that each element of $\Lambda(V)$ is a linear combination of elements

$$
v_{i_{1}} \wedge \cdots \wedge v_{i_{k}} \quad \text { for } i_{1}<\cdots<i_{k}
$$

and that the relations cannot be used further. Thus we have shown:

Proposition 13.5 (i) $\Lambda(V)=\Lambda^{0} V \oplus \Lambda^{1} V \oplus \cdots \oplus \Lambda^{n} V$
(ii) $\operatorname{dim} \Lambda^{k} V=\binom{n}{k}$
(iii) $\operatorname{dim} \Lambda(V)=2^{n}$
(iv) The elements $v_{i_{1}} \wedge \cdots \wedge v_{i_{k}}$ for subsets $\left\{i_{1}, \ldots, i_{k}\right\} \subset\{1, \ldots, n\}$ with $i_{1}<\cdots<i_{k}$ form a basis of $\Lambda(V)$.

We now show that $\Lambda(V)$ has the structure of an $L$-module. We recall that $T(V)$ is an $L$-module and that its ideal $J$ is generated by $v \otimes v$ for all $v \in V$. For $x \in L$ we have

$$
x(v \otimes v)=x v \otimes v+v \otimes x v
$$

Since the right-hand side lies in $J$ we see that $J$ is a submodule of $T(V)$. Thus $\Lambda(V)=T(V) / J$ can be made into an $L$-module in the natural way. Each exterior power $\Lambda^{k} V$ is a submodule.

Proposition 13.6 Let $V$ be a finite dimensional module for the simple Lie algebra $L$. Then the weights of $\Lambda^{k} V$ are all sums of $k$ distinct weights of $V$.

Proof. Let $H$ be a Cartan subalgebra of $L$. We consider $V$ as an $H$-module. $V$ is a direct sum of 1 -dimensional $H$-submodules. Let $v_{1}, \ldots, v_{n}$ be a basis of $V$ adapted to this decomposition. Let $\lambda_{1}, \ldots, \lambda_{n} \in H^{*}$ be the corresponding weights. Then

$$
x v_{i}=\lambda_{i}(x) v_{i} \quad \text { for } x \in H
$$

Now $\Lambda^{k} V$ has basis $v_{i_{1}} \wedge \cdots \wedge v_{i_{k}}$ for all $i_{1}<\cdots<i_{k}$. We have

$$
\begin{aligned}
x\left(v_{i_{1}} \wedge \cdots \wedge v_{i_{k}}\right) & =\sum_{r=1}^{k} v_{i_{1}} \wedge \cdots \wedge x v_{i_{r}} \wedge \cdots \wedge v_{i_{k}} \\
& =\left(\lambda_{i_{1}}(x)+\cdots+\lambda_{i_{k}}(x)\right) v_{i_{1}} \wedge \cdots \wedge v_{i_{k}} \quad x \in H
\end{aligned}
$$

Thus $v_{i_{1}} \wedge \cdots \wedge v_{i_{k}}$ is a weight vector with weight $\lambda_{i_{1}}+\cdots+\lambda_{i_{k}}$. Thus the weights of $\Lambda^{k} V$ are sums of $k$ distinct weights of $V$.

Theorem 13.7 Let $V$ be the natural module for the simple Lie algebra $A_{l}$. Then the fundamental modules for $A_{l}$ are

$$
\Lambda^{1} V, \Lambda^{2} V, \ldots, \Lambda^{l} V
$$

Proof. We have seen in Proposition 13.3 that $\Lambda^{1} V=V$ is the fundamental module with highest weight $\omega_{1}$. The weights of $V$ are the maps $\mu_{1}, \ldots, \mu_{l+1}$ given by

$$
\mu_{i}:\left(\begin{array}{cccc}
\lambda_{1} & & & \\
& \lambda_{2} & & \\
& & \ddots & \\
& & & \lambda_{l+1}
\end{array}\right) \rightarrow \lambda_{i}
$$

Since $\mu_{i}-\mu_{i+1}=\alpha_{i}$ we have $\mu_{i} \succ \mu_{i+1}$. Thus the weights are ordered by

$$
\mu_{1} \succ \cdots \succ \mu_{l+1}
$$

Now $\mu_{1}=\omega_{1}, \mu_{i}-\mu_{i+1}=\alpha_{i}$ and $\alpha_{i}=\sum_{j} A_{j i} \omega_{j}$ by Proposition 10.17. Hence

$$
\begin{aligned}
& \alpha_{1}=2 \omega_{1}-\omega_{2} \\
& \alpha_{i}=-\omega_{i-1}+2 \omega_{i}-\omega_{i+1} \quad \text { for } 2 \leq i \leq l-1 \\
& \alpha_{l}=-\omega_{l-1}+2 \omega_{l} .
\end{aligned}
$$

It follows that

$$
\begin{aligned}
& \mu_{1}=\omega_{1} \\
& \mu_{2}=-\omega_{1}+\omega_{2} \\
& \mu_{3}=-\omega_{2}+\omega_{3} \\
& \vdots \\
& \mu_{l}=-\omega_{l-1}+\omega_{l} \\
& \mu_{l+1}=-\omega_{l} .
\end{aligned}
$$

By Proposition 13.6 the highest weight of $\Lambda^{k} V$ is $\mu_{1}+\cdots+\mu_{k}=\omega_{k}$, for $1 \leq k \leq l$. Thus $\Lambda^{k} V$ contains the irreducible module $L\left(\omega_{k}\right)$ as one of its irreducible direct summands. However,

$$
\operatorname{dim} L\left(\omega_{k}\right)=\operatorname{dim} \Lambda^{k} V=\binom{l+1}{k}
$$

by Proposition 13.2. Hence $L\left(\omega_{k}\right)=\Lambda^{k} V$.

### 13.4 Fundamental modules for $B_{l}$ and $D_{l}$

The fundamental weights $\omega_{1}, \ldots, \omega_{l}$ for a simple Lie algebra of type $B_{l}$ or $D_{l}$ will be numbered according to the labelling of the Dynkin diagrams:


We again use Theorem 13.1 to calculate $\operatorname{dim} L\left(\omega_{j}\right)$.
We suppose first that we have an algebra of type $B_{l}$. We know from Section 8.3 that the roots have the following form. Let

$$
h=\left(\begin{array}{lllllll}
0 & & & & & & \\
& \lambda_{1} & & & & & \\
& & \ddots & & & & \\
& & & \lambda_{l} & & & \\
& & & & -\lambda_{1} & & \\
& & & & & \ddots & \\
& & & & & & \\
& & & & & & \\
& & & & & \\
&
\end{array}\right)
$$

Then the fundamental roots are

$$
\begin{aligned}
& \alpha_{i}(h)=\lambda_{i}-\lambda_{i+1} \quad \text { for } 1 \leq i \leq l-1 \\
& \alpha_{l}(h)=\lambda_{l}
\end{aligned}
$$

The full set of positive roots is given by

$$
\begin{array}{ll}
h \rightarrow \lambda_{i}-\lambda_{j} & \text { for } i<j \\
h \rightarrow \lambda_{i}+\lambda_{j} & \text { for } i<j \\
h \rightarrow \lambda_{i} &
\end{array}
$$

where $i, j \in\{1, \ldots, l\}$. These positive roots can be expressed as combinations of fundamental roots as follows:

$$
\begin{aligned}
& \alpha_{i}+\cdots+\alpha_{j-1} \quad \text { for } i<j \\
& \alpha_{i}+\cdots+\alpha_{j-1}+2 \alpha_{j}+\cdots+2 \alpha_{l} \quad \text { for } i<j \\
& \alpha_{i}+\cdots+\alpha_{l}
\end{aligned}
$$

The first two families are long roots and the third family are short roots. Thus the weights $w_{i}$ are given by

$$
w_{1}=\cdots=w_{l-1}=2 \quad w_{l}=1
$$

According to Theorem 13.1 we have

$$
\operatorname{dim} L\left(\omega_{j}\right)=\prod_{\alpha \in \Phi^{+}} d_{\alpha}
$$

where $\alpha=\sum k_{i} \alpha_{i}$ and

$$
d_{\alpha}=\frac{\sum_{i=1}^{l} k_{i} w_{i}+k_{j} w_{j}}{\sum_{i=1}^{l} k_{i} w_{i}}
$$

We have $d_{\alpha}=1$ if $\alpha$ does not involve $\alpha_{j}$. We first suppose $j \in\{1, \ldots, l-1\}$. Then the positive roots involving $j$ are:

$$
\begin{aligned}
& \alpha_{i}+\cdots+\alpha_{j}+\cdots+\alpha_{k} \quad 1 \leq i \leq j, j \leq k \leq l-1 \\
& \alpha_{i}+\cdots+\alpha_{j}+\cdots+\alpha_{l} \quad 1 \leq i \leq j \\
& \alpha_{i}+\cdots+\alpha_{j}+\cdots+\alpha_{k-1}+2 \alpha_{k}+\cdots+2 \alpha_{l} \quad 1 \leq i \leq j, \quad j \leq k-1<l \\
& \alpha_{i}+\cdots+\alpha_{k-1}+2 \alpha_{k}+\cdots+2 \alpha_{j}+\cdots+2 \alpha_{l} \quad 1 \leq i<k, \quad k \leq j \leq l-1
\end{aligned}
$$

The values of $d_{\alpha}$ in these four cases are

$$
\frac{k-i+2}{k-i+1}, \quad \frac{2 l-2 i+3}{2 l-2 i+1}, \quad \frac{2 l-k-i+2}{2 l-k-i+1}, \quad \frac{2 l-k-i+3}{2 l-k-i+1}
$$

respectively. The product of all possible $d_{\alpha}$ in these four cases is

$$
\begin{aligned}
& \frac{(j+1)(j+2) \cdots l}{1 \cdot 2 \cdots l-j}, \quad \frac{2 l+1}{2 l-2 j+1}, \quad \frac{(2 l-j)(2 l-j-1) \cdots(l+1)}{(2 l-2 j)(2 l-2 j-1) \cdots(l+1-j)}, \\
& \frac{2 l(2 l-1)(2 l-2) \cdots(2 l+2-j)}{(2 l-j)(2 l-j-1) \cdots(2 l-2 j+3)(2 l-2 j+2)}
\end{aligned}
$$

respectively. Finally the total product $\prod_{\alpha \in \Phi^{+}} d_{\alpha}$ is $\binom{2 l+1}{j}$. We now take $j=l$. Then the positive roots involving $l$ are

$$
\begin{aligned}
& \alpha_{i}+\cdots+\alpha_{l} \quad 1 \leq i \leq l \\
& \alpha_{i}+\cdots+\alpha_{k-1}+2 \alpha_{k}+\cdots+2 \alpha_{l} \quad 1 \leq i<k \leq l
\end{aligned}
$$

The values of $d_{\alpha}$ in these cases are $\frac{2 l-2 i+2}{2 l-2 i+1}, \quad \frac{2 l-i-j+2}{2 l-i-j+1}$ respectively. The product of all possible $d_{\alpha}$ in the two cases is

$$
\frac{2 l(2 l-2)(2 l-4) \cdots 2}{(2 l-1)(2 l-3) \cdots 3 \cdot 1}, \quad \frac{(2 l-1)(2 l-2) \cdots(l+1)}{(2 l-2)(2 l-4) \cdots 2}
$$

respectively. Finally the total product $\prod_{\alpha \in \Phi^{+}} d_{\alpha}$ is $2^{l}$.
Thus we have shown:

Proposition 13.8 The dimensions of the fundamental modules for the simple Lie algebra of type $B_{l}$ are


The dimensions of the modules $L\left(\omega_{j}\right)$ for $1 \leq j \leq l-1$ suggest that these modules are exterior powers of the $(2 l+1)$-dimensional natural module. This is indeed the case.

Theorem 13.9 Let $V$ be the $(2 l+1)$-dimensional natural module for the simple Lie algebra $B_{l}$ (described in Section 8.3). Then the fundamental module $L\left(\omega_{j}\right)$ is isomorphic to $\Lambda^{j} V$ for $1 \leq j \leq l-1$.

Proof. Let

$$
h=\left(\begin{array}{lllllll}
0 & & & & & & \\
& \lambda_{1} & & & & & \\
& & \ddots & & & & \\
& & & \lambda_{l} & & & \\
& & & & -\lambda_{1} & & \\
& & & & & \ddots & \\
& & & & & & -\lambda_{l}
\end{array}\right)
$$

Then the weights of $V$ are $0, \mu_{1}, \ldots, \mu_{l},-\mu_{1}, \ldots,-\mu_{l}$ where $\mu_{i}(h)=\lambda_{i}$. Since $\mu_{i}-\mu_{i+1}=\alpha_{i}$ for $1 \leq i \leq l-1$ and $\mu_{l}=\alpha_{l}$ we have

$$
\mu_{1} \succ \mu_{2} \succ \cdots \succ \mu_{l} \succ 0
$$

Thus the highest weight of $\Lambda^{j} V$ for $1 \leq j \leq l$ is $\mu_{1}+\mu_{2}+\cdots+\mu_{j}$. Expressing the $\mu \mathrm{s}$ in terms of the $\alpha \mathrm{s}$ gives

$$
\begin{aligned}
& \mu_{1}=\alpha_{1}+\cdots+\alpha_{l} \\
& \mu_{2}=\alpha_{2}+\cdots+\alpha_{l}
\end{aligned}
$$

$$
\mu_{l}=\alpha_{l}
$$

We also have $\alpha_{i}=\sum A_{j i} \omega_{j}$, which in type $B_{l}$ gives

$$
\begin{aligned}
& \alpha_{1}=2 \omega_{1}-\omega_{2} \\
& \alpha_{i}=-\omega_{i-1}+2 \omega_{i}-\omega_{i+1} \quad 2 \leq i \leq l-2 \\
& \alpha_{l-1}=-\omega_{l-2}+2 \omega_{l-1}-2 \omega_{l} \\
& \alpha_{l}=\quad-\omega_{l-1}+2 \omega_{l} .
\end{aligned}
$$

It follows that

$$
\begin{aligned}
& \mu_{1}=\omega_{1} \\
& \mu_{2}=-\omega_{1}+\omega_{2} \\
& \vdots \\
& \mu_{l-1}=-\omega_{l-2}+\omega_{l-1} \\
& \mu_{l}=-\omega_{l-1}+2 \omega_{l} .
\end{aligned}
$$

Hence $\mu_{1}+\cdots+\mu_{j}=\omega_{j}$ for $1 \leq j \leq l-1$

$$
\mu_{1}+\cdots+\mu_{l}=2 \omega_{l}
$$

Thus the highest weight of $\Lambda^{j} V$ is $\omega_{j}$ for $1 \leq j \leq l-1$. Since

$$
\operatorname{dim} L\left(\omega_{j}\right)=\operatorname{dim} \Lambda^{j} V=\binom{2 l+1}{j} \quad j \leq l-1
$$

we deduce that $L\left(\omega_{j}\right)$ is isomorphic to $\Lambda^{j} V$. This argument fails when $j=l$ since the highest weight of $\Lambda^{l} V$ is $2 \omega_{l}$ rather than $\omega_{l}$. We shall see subsequently how to find the remaining fundamental module $L\left(\omega_{l}\right)$.

We now consider the simple Lie algebra of type $D_{l}$. This algebra was described in Section 8.2. Its roots have the following form. Let

$$
h=\left(\begin{array}{llllll}
\lambda_{1} & & & & & \\
& \ddots & & & & \\
& & & & & \\
l & \\
& & & -\lambda_{1} & & \\
& & & & \ddots & \\
& & & & & \\
& & & & & -\lambda_{l}
\end{array}\right)
$$

Then the fundamental roots are

$$
\begin{aligned}
& \alpha_{i}(h)=\lambda_{i}-\lambda_{i+1} \quad \text { for } 1 \leq i \leq l-1 \\
& \alpha_{l}(h)=\lambda_{l-1}+\lambda_{l} .
\end{aligned}
$$

The full set of positive roots is given by

$$
\begin{array}{ll}
h \rightarrow \lambda_{i}-\lambda_{j} & i<j \\
h \rightarrow \lambda_{i}+\lambda_{j} & i<j
\end{array}
$$

These are expressed as combinations of the fundamental roots by

$$
\begin{aligned}
& \alpha_{i}+\cdots+\alpha_{j-1} \quad \text { for } 1 \leq i<j \leq l \\
& \alpha_{i}+\cdots+\alpha_{j-1}+2 \alpha_{j}+\cdots+2 \alpha_{l-2}+\alpha_{l-1}+\alpha_{l} \quad \text { for } 1 \leq i<j \leq l-1 \\
& \alpha_{i}+\cdots+\alpha_{l-2}+\alpha_{l} \quad \text { for } 1 \leq i \leq l-2
\end{aligned}
$$

We take a fixed $j$ with $1 \leq j \leq l-2$ and consider $\operatorname{dim} L\left(\omega_{j}\right)$. By Theorem 13.1 this is given by

$$
\operatorname{dim} L\left(\omega_{j}\right)=\prod_{\alpha \in \Phi^{+}} d_{\alpha}
$$

where $\alpha=\sum k_{i} \alpha_{i}$ and

$$
d_{\alpha}=\frac{\sum_{i=1}^{l} k_{i}+k_{j}}{\sum_{i=1}^{l} k_{i}}
$$

(All weights $w_{i}$ are equal to 1 in type $D_{l}$.) As usual $d_{\alpha}=1$ if $\alpha$ does not involve $\alpha_{j}$. The positive roots involving $\alpha_{j}$ are
$\alpha_{i}+\cdots+\alpha_{j}+\cdots+\alpha_{k} \quad 1 \leq i \leq j, j \leq k \leq l-1$
$\alpha_{i}+\cdots+\alpha_{j}+\cdots+\alpha_{k}+2 \alpha_{k+1}+\cdots+2 \alpha_{l-2}+\alpha_{l-1}+\alpha_{l} \quad 1 \leq i \leq j, j \leq k \leq l-2$
$\alpha_{i}+\cdots+\alpha_{k}+2 \alpha_{k+1}+\cdots+2 \alpha_{j}+\cdots+2 \alpha_{l-2}+\alpha_{l-1}+\alpha_{l} \quad 1 \leq i \leq k<j$
$\alpha_{i}+\cdots+\alpha_{j}+\cdots+\alpha_{l-2}+\alpha_{l} \quad 1 \leq i \leq j$.
The values of $d_{\alpha}$ in these four cases are

$$
\frac{k-i+2}{k-i+1}, \quad \frac{2 l-i-k}{2 l-i-k-1}, \quad \frac{2 l-i-k+1}{2 l-i-k-1}, \quad \frac{l-i+1}{l-i}
$$

respectively. The product of all possible $d_{\alpha}$ in these four cases is

$$
\begin{aligned}
& \frac{(j+1)(j+2) \cdots l}{1 \cdot 2 \cdots(l-j)}, \quad \frac{(2 l-j-1)(2 l-j-2) \cdots(l+1)}{(2 l-2 j-1)(2 l-2 j-2) \cdots(l+1-j)}, \\
& \frac{(2 l-1)(2 l-2) \cdots(2 l-j+1)}{(2 l-j-1)(2 l-j-2) \cdots(2 l-2 j+1)}, \quad \frac{l}{l-j}
\end{aligned}
$$

respectively. Finally the total product is $\binom{2 l}{j}$.
We next suppose $j=l-1$. The positive roots involving $\alpha_{l-1}$ are

$$
\begin{aligned}
& \alpha_{i}+\cdots+\alpha_{l-1} \quad 1 \leq i \leq l-1 \\
& \alpha_{i}+\cdots+\alpha_{k-1}+2 \alpha_{k}+\cdots+2 \alpha_{l-2}+\alpha_{l-1}+\alpha_{l} \quad 1 \leq i<k \leq l-1
\end{aligned}
$$

The values of $d_{\alpha}$ in these two cases are $\frac{l-i+1}{l-i}, \frac{2 l-i-k+1}{2 l-i-k}$ respectively. The product of all possible $d_{\alpha}$ in those two cases is $l, 2^{l-1} / l$ respectively, and so the total product is $2^{l-1}$.

Finally suppose $j=l$. The positive roots involving $\alpha_{l}$ are

$$
\begin{aligned}
& \alpha_{l} \\
& \alpha_{i}+\cdots+\alpha_{l-2}+\alpha_{l} \quad 1 \leq i \leq l-2 \\
& \alpha_{i}+\cdots+\alpha_{k-1}+2 \alpha_{k}+\cdots+2 \alpha_{l-2}+\alpha_{l-1}+\alpha_{l} \quad 1 \leq i<k \leq l-1
\end{aligned}
$$

The values of $d_{\alpha}$ in these three cases are $2, \quad \frac{l-i+1}{l-i}, \frac{2 l-i-k+1}{2 l-i-k}$ respectively. The product of all possible $d_{\alpha}$ in these three cases is $2, l / 2,2^{l-1} / l$ respectively. Thus the total product is $2^{l-1}$.

Thus we have shown
Proposition 13.10 The dimensions of the fundamental modules for the simple Lie algebra of type $D_{l}$ are


Again the dimensions of these modules for $1 \leq j \leq l-2$ suggest that they are given by exterior powers of the natural module.

Theorem 13.11 Let $V$ be the 2l-dimensional natural module for the simple Lie algebra $D_{l}$ (described in Section 8.2). Then the fundamental module $L\left(\omega_{j}\right)$ is isomorphic to $\Lambda^{j} V$ for $1 \leq j \leq l-2$.

## Proof. Let

$$
h=\left(\begin{array}{llllll}
\lambda_{1} & & & & & \\
& \ddots & & & & \\
& & \lambda_{l} & & & \\
& & & -\lambda_{1} & & \\
& & & & \ddots & \\
& & & & & -\lambda_{l}
\end{array}\right)
$$

Then the weights of $V$ are $\mu_{1}, \ldots, \mu_{l},-\mu_{1}, \ldots,-\mu_{l}$ where $\mu_{i}(h)=\lambda_{i}$. We have

$$
\begin{aligned}
& \mu_{1}-\mu_{2}=\alpha_{1} \\
& \vdots \\
& \mu_{l-1}-\mu_{l}=\alpha_{l-1} \\
& \mu_{l-1}+\mu_{l}=\alpha_{l}
\end{aligned}
$$

and

$$
\begin{aligned}
& \alpha_{1}=2 \omega_{1}-\omega_{2} \\
& \alpha_{2}=-\omega_{1}+2 \omega_{2}-\omega_{3} \\
& \vdots \\
& \alpha_{l-3}=-\omega_{l-4}+2 \omega_{l-3}-\omega_{l-2}
\end{aligned}
$$

$$
\begin{aligned}
& \alpha_{l-2}=-\omega_{l-3}+2 \omega_{l-2}-\omega_{l-1}-\omega_{l} \\
& \alpha_{l-1}=-\omega_{l-2}+2 \omega_{l-1} \\
& \alpha_{l}=-\omega_{l-2}+2 \omega_{l}
\end{aligned}
$$

using the Cartan matrix of type $D_{l}$. It follows that

$$
\begin{aligned}
& \mu_{1}=\omega_{1} \\
& \mu_{2}=-\omega_{1}+\omega_{2} \\
& \vdots \\
& \mu_{l-2}=-\omega_{l-3}+\omega_{l-2} \\
& \mu_{l-1}=-\omega_{l-2}+\omega_{l-1}+\omega_{l} \\
& \mu_{l}=-\omega_{l-1}+\omega_{l} .
\end{aligned}
$$

Since $\mu_{1} \succ \mu_{2} \succ \cdots \succ \mu_{l-1} \succ \mu_{l}$ the highest weight of $\Lambda^{j} V$ for $1 \leq j \leq l-2$ is $\mu_{1}+\cdots+\mu_{j}$. Also we have $\mu_{1}+\cdots+\mu_{j}=\omega_{j}$ for $j \leq l-2$. Thus the highest weight of $\Lambda^{j} V$ is $\omega_{j}$ when $j \leq l-2$. Since

$$
\operatorname{dim} L\left(\omega_{j}\right)=\operatorname{dim} \Lambda^{j} V=\binom{2 l}{j}, \quad j \leq l-2
$$

we deduce that $L\left(\omega_{j}\right)$ is isomorphic to $\Lambda^{j} V$ for $j \leq l-2$.

### 13.5 Clifford algebras and spin modules

There remain one fundamental module for $B_{l}$ of dimension $2^{l}$ and two fundamental modules for $D_{l}$ of dimension $2^{l-1}$ which cannot be obtained as exterior powers of the natural module. These are called spin modules and give rise to spin representations of $B_{l}$ and $D_{l}$. We shall now show how these modules may be obtained in terms of the Clifford algebra.

Let $V$ be a vector space of dimension $n$ over $\mathbb{C}$ and suppose we are given a symmetric bilinear map $V \times V \rightarrow \mathbb{C}$ under which the pair $v, v^{\prime}$ maps to $\left(v, v^{\prime}\right) \in \mathbb{C}$. Thus we have

$$
\left(v^{\prime}, v\right)=\left(v, v^{\prime}\right)
$$

Let $T(V)$ be the tensor algebra of $V$ and $J$ be the two-sided ideal of $T(V)$ generated by elements

$$
v \otimes v-(v, v) 1 \quad \text { for all } v \in V
$$

Since

$$
\begin{aligned}
\left(v+v^{\prime}\right) \otimes\left(v+v^{\prime}\right)-\left(v+v^{\prime}, v+v^{\prime}\right) 1= & (v \otimes v-(v, v) 1)+\left(v^{\prime} \otimes v^{\prime}-\left(v^{\prime}, v^{\prime}\right) 1\right) \\
& +\left(v \otimes v^{\prime}+v^{\prime} \otimes v-2\left(v, v^{\prime}\right) 1\right)
\end{aligned}
$$

we see that

$$
v \otimes v^{\prime}+v^{\prime} \otimes v-2\left(v, v^{\prime}\right) 1 \in J \quad \text { for all } v, v^{\prime} \in V
$$

Let $C(V)=T(V) / J$. Then $C(V)$ is an associative algebra called the Clifford algebra of $V$.

Now let $v_{1}, \ldots, v_{n}$ be a basis of $V$. Then the elements

$$
\begin{aligned}
& v_{i} \otimes v_{i}-\left(v_{i}, v_{i}\right) 1 \\
& v_{i} \otimes v_{j}+v_{j} \otimes v_{i}-2\left(v_{i}, v_{j}\right) 1 \quad i<j
\end{aligned}
$$

lie in $J$ and it is evident that these elements generate $J$ as a 2 -sided ideal. We observe also that

$$
(\mathbb{C} 1 \oplus V) \cap J=\left(T^{0}(V) \oplus T^{1}(V)\right) \cap J=O
$$

and so the natural map $T(V) \rightarrow C(V)$ is injective when restricted to $\mathbb{C} 1 \oplus V$. We shall regard $\mathbb{C} 1 \oplus V$ as a subspace of $C(V)$. Thus $C(V)$ is generated, as associative algebra with 1 , by elements $v_{1}, \ldots, v_{n}$ subject to relations

$$
\begin{aligned}
& v_{i} v_{i}=\left(v_{i}, v_{i}\right) 1 \\
& v_{j} v_{i}=-v_{i} v_{j}+2\left(v_{i}, v_{j}\right) 1 \quad i<j
\end{aligned}
$$

By using these relations any polynomial in $v_{1}, \ldots, v_{n}$ can be written as a polynomial in which each monomial has form $v_{i_{1}} v_{i_{2}} \ldots v_{i_{k}}$ where $i_{1}<i_{2}<\cdots<i_{k}$ and $0 \leq k \leq n$. Moreover an element of $C(V)$ in this standard form cannot be simplified further by the use of the above relations. Thus we have shown

Proposition 13.12 (i) $\operatorname{dim} C(V)=2^{n}$.
(ii) The elements $v_{i_{1}} v_{i_{2}} \ldots v_{i_{k}}$ for $i_{1}<i_{2}<\cdots<i_{k}$ with $0 \leq k \leq n$ form a basis for $C(V)$. (The empty product is 1 .)

We note that all generators of $J$ lie in $T^{0} V \oplus T^{2} V$. We define $T(V)^{+}, T(V)^{-}$by

$$
\begin{aligned}
& T(V)^{+}=\bigoplus_{i \text { even }} T^{i} V \\
& T(V)^{-}=\bigoplus_{i \text { odd }} T^{i} V
\end{aligned}
$$

Then we have

$$
\begin{aligned}
& T(V)=T(V)^{+} \oplus T(V)^{-} \\
& J=\left(J \cap T(V)^{+}\right) \oplus\left(J \cap T(V)^{-}\right)
\end{aligned}
$$

This follows from the fact that $J$ is generated by elements of $T(V)^{+}$. Hence

$$
C(V) \cong \frac{T(V)^{+}}{J \cap T(V)^{+}} \oplus \frac{T(V)^{-}}{J \cap T(V)^{-}}
$$

We write $C(V)^{+}=\frac{T(V)^{+}}{J \cap T\left(V^{+}\right)}$and $C(V)^{-}=\frac{T(V)^{-}}{J \cap T(V)^{-}}$. Then

$$
C(V)=C(V)^{+} \oplus C(V)^{-}
$$

In terms of our basis for $C(V), C(V)^{+}$has basis $v_{i_{1}} v_{i_{2}} \ldots v_{i_{k}}$ for $i_{1}<i_{2}<\cdots<$ $i_{k}$ with $k$ even and $C(V)^{-}$has basis consisting of these elements with $k$ odd. Thus

$$
\operatorname{dim} C(V)^{+}=\operatorname{dim} C(V)^{-}=2^{n-1}
$$

Now the associative algebra $C(V)$ can be made into a Lie algebra $[C(V)]$ in the usual way by defining $[x y]=x y-y x$. Let $L$ be the subspace of $[C(V)]$ spanned by the elements $\left[v v^{\prime}\right]$ for all $v, v^{\prime} \in V$. Then $L$ can be spanned by elements $\left[v_{i} v_{j}\right]$ for $i<j$, and since

$$
\left[v_{i} v_{j}\right]=2 v_{i} v_{j}-2\left(v_{i}, v_{j}\right) 1
$$

these elements are linearly independent. Thus $\operatorname{dim} L=n(n-1) / 2$. We shall show that $L$ is a Lie subalgebra of $[C(V)]$.

Lemma 13.13 (i) Let $x, y, z \in V$. Then $[[x y] z]=4(y, z) x-4(x, z) y$.
(ii) Let $x, y, z, w \in V$. Then

$$
[[x y],[z w]]=4(y, z)[x w]-4(y, w)[x z]+4(x, w)[y z]-4(x, z)[y w]
$$

Proof. (i) $[[x y] z]=(x y-y x) z-z(x y-y x)$

$$
\begin{aligned}
= & x y z-y x z-z x y+z y x \\
= & -x z y+2(y, z) x+y z x-2(x, z) y \\
& +x z y-2(x, z) y-y z x+2(y, z) x \\
= & 4(y, z) x-4(x, z) y .
\end{aligned}
$$

(ii) $[[x y],[z w]]=[x y] z w-[x y] w z-z w[x y]+w z[x y]$

$$
\begin{aligned}
& =[[x y] z] w+z[x y] w-[[x y] w] z-w[x y] z-z w[x y]+w z[x y] \\
& =[[x y] z] w-[[x y] w] z+z[[x y] w]-w[[x y] z] \\
& =[[[x y] z] w]-[[[x y] w] z] \\
& =[4(y, z) x-4(x, z) y, w]-[4(y, w) x-4(x, w) y, z] \\
& =4(y, z)[x w]-4(x, z)[y w]-4(y, w)[x z]+4(x, w)[y z] .
\end{aligned}
$$

Corollary 13.14 L is a Lie subalgebra of $[C(V)]$.

Now $C(V)$ is a $[C(V)]$-module giving the adjoint representation so is in particular an $L$-module. Lemma 13.13 (i) shows that its subspace $V$ is an $L$-submodule.

Proposition 13.15 Suppose the symmetrix scalar product on $V$ is nondegenerate. Then $V$ is a faithful L-module.

Proof. Let $x \in L$ and suppose $[x v]=0$ for all $v \in V$. We must show $x=0$. Let $x=\sum_{i<j} c_{i j}\left[v_{i} v_{j}\right]$. We may define a skew-symmetrix $n \times n$ matrix $C=\left(c_{i j}\right)$ by $c_{i i}=0$ and $c_{j i}=-c_{i j}$ for $i<j$. We have

$$
\sum_{i<j} c_{i j}\left[\left[v_{i} v_{j}\right] v\right]=0
$$

By Lemma 13.13 (i) we have

$$
4 \sum_{i<j} c_{i j}\left(\left(v_{j}, v\right) v_{i}-\left(v_{i}, v\right) v_{j}\right)=0 .
$$

The coefficient of $v_{i}$ in this expression is

$$
4\left(\sum_{\substack{j \\ j>i}} c_{i j}\left(v_{j}, v\right)-\sum_{\substack{j \\ j<i}} c_{j i}\left(v_{j}, v\right)\right)=4 \sum_{j=1}^{n} c_{i j}\left(v_{j}, v\right) .
$$

It follows that

$$
\sum_{j=1}^{n} c_{i j}\left(v_{j}, v\right)=0
$$

for all $i$ and all $v \in V$. Let $\left(v_{j}, v_{k}\right)=m_{j k}$. Then we have

$$
\sum_{j=1}^{n} c_{i j} m_{j k}=0 \quad \text { for all } i, k
$$

that is $C M=O$ where $M=\left(m_{j k}\right)$. If the scalar product on $V$ is non-degenerate then $M$ is a non-singular matrix. Then $C=O$ and so $x=0$.

Lemma 13.16 Let $x \in L$ and $v, v^{\prime} \in V$. Then

$$
\left([x v], v^{\prime}\right)+\left(v,\left[x v^{\prime}\right]\right)=0 .
$$

Proof. It is sufficient to prove this when $x=[y z]$ for $y, z \in V$. Now

$$
\begin{aligned}
& \left([[y z] v], v^{\prime}\right)=4(z, v)\left(y, v^{\prime}\right)-4(y, v)\left(z, v^{\prime}\right) \\
& \left(v,\left[[y z] v^{\prime}\right]\right)=4\left(z, v^{\prime}\right)(v, y)-4\left(y, v^{\prime}\right)(v, z)
\end{aligned}
$$

by Lemma 13.13 (i). The result follows.

Thus we have a Lie algebra $L$ of dimension $n(n-1) / 2$, an $L$-module $V$ of dimension $n$, and a symmetric bilinear scalar product on $V$ invariant under $L$ in the sense of Lemma 13.16.

We now consider some special cases of the above situation. First let $V$ be a vector space with $\operatorname{dim} V=2 l+1$ and let $v_{0}, v_{1}, \ldots, v_{l}, v_{-1}, \ldots, v_{-l}$ be a basis of $V$. Consider the symmetric bilinear scalar product $V \times V \rightarrow \mathbb{C}$ determined by

$$
\begin{aligned}
& \left(v_{0}, v_{0}\right)=2 \\
& \left(v_{i}, v_{-i}\right)=1 \quad i=1, \ldots, l
\end{aligned}
$$

and all other scalar products of basis elements 0 . The matrix of this scalar product is

$$
\left(\begin{array}{ccccc}
2 & 0 & \cdots & \cdots & 0 \\
0 & O & & I_{l} \\
\vdots & & & \\
0 & I_{l} & & & O
\end{array}\right)
$$

The condition

$$
\left([x v], v^{\prime}\right)+\left(v,\left[x v^{\prime}\right]\right)=0
$$

of Lemma 13.16 tells us that if the element $x \in L$ is represented by the matrix $X$ on $V$ then

$$
X^{\mathrm{t}} M+M X=O
$$

Now $L$ has dimension $l(2 l+1)$ and, by Proposition 13.15 , acts faithfully on $V$. However, the Lie algebra of all $(2 l+1) \times(2 l+1)$ matrices $X$ satisfying $X^{\mathrm{t}} M+M X=O$ is the simple Lie algebra $B_{l}$ (cf. Section 8.3) and has dimension $l(2 l+1)$. Since $L$ is contained in this set of matrixes we must have $L=B_{l}$.

We aim to find the spin module for $B_{l}$ inside the Clifford algebra $C(V)$. Recall that $V$ has basis

$$
v_{0}, v_{1}, \ldots, v_{l}, v_{-1}, \ldots, v_{-l}
$$

We define $u_{i}=v_{0} v_{i}$ and $u_{-i}=v_{0} v_{-i}$ for $i=1, \ldots, l$. Let $U$ be the subspace of $C(V)$ spanned by elements

$$
u_{-j_{1}} u_{-j_{2}} \ldots u_{-j_{t}} u_{1} u_{2} \ldots u_{l}
$$

for all subsets $j_{1}<j_{2}<\cdots<j_{t}$ of $\{1, \ldots, l\}$ with $0 \leq t \leq l$. Bearing in mind the natural basis of $C(V)$ we see that these elements are linearly independent. Thus $\operatorname{dim} U=2^{l}$. We have $U \subset C(V)^{+}$.

Lemma 13.17 $C(V)^{+} U \subset U$. Thus $U$ is a left ideal of $C(V)^{+}$.

Proof. We first observe that $C(V)^{+}$is generated by the elements $u_{i}$ and $u_{-i}$ for $i=1, \ldots, l$. This follows from the fact that $v_{0}$ anticommutes with $v_{i}$ and $v_{-i}$ for all $i=1, \ldots, l$.

We note next that $u_{i}^{2}=0$ and $u_{-i}^{2}=0$ for $i=1, \ldots, l$, that, $u_{i}, u_{j}$ anticommute and $u_{-i}, u_{-j}$ anticommute when $i \neq j$ and both lie in $\{1, \ldots, l\}$, that $u_{i}, u_{-j}$ anticommute for all $i, j \in\{1, \ldots, l\}$, and that

$$
u_{i} u_{-i}+u_{-i} u_{i}=-4 \cdot 1 .
$$

For $u_{i} u_{-i}+u_{-i} u_{i}=v_{0} v_{i} v_{0} v_{-i}+v_{0} v_{-i} v_{0} v_{i}=-v_{0}^{2}\left(v_{i} v_{-i}+v_{-i} v_{i}\right)=-2 \cdot 2\left(v_{i}, v_{-i}\right) 1$ $=-4 \cdot 1$.

It follows from these relations that

$$
\begin{aligned}
& u_{i} \cdot u_{-j_{1}} \ldots u_{-j_{t}} u_{1} \ldots u_{l} \\
& \quad= \begin{cases} \pm 4 u_{-j_{1}} \ldots \hat{u}_{-i} \ldots u_{-j_{t}} u_{1} \ldots u_{l} & \text { if } i \in\left\{j_{1}, \ldots, j_{t}\right\} \\
0 & \text { otherwise }\end{cases}
\end{aligned}
$$

where $\hat{u}_{-i}$ means the term $u_{-i}$ is omitted.
$u_{-i} \cdot u_{-j_{1}} \ldots u_{-j_{t}} u_{1} \ldots u_{l}= \begin{cases}0 & \text { if } i \in\left\{j_{1}, \ldots, j_{t}\right\} \\ \pm u_{-j_{1}} \ldots u_{-i} \ldots u_{-j_{t}} u_{1} \ldots u_{l} & \text { if } i \notin\left\{j_{1}, \ldots, j_{t}\right\} .\end{cases}$
This shows that $u_{i} U \subset U$ and $u_{-i} U \subset U$, so $C(V)^{+} U \subset U$.

This lemma shows that we may regard $U$ as a $C(V)^{+}$-module under left multiplication. $U$ is therefore a $\left[C(V)^{+}\right]$-module under the same action. Since $L$ is a Lie subalgebra of $\left[C(V)^{+}\right]$we may regard $U$ as an $L$-module under left multiplication.

Warning note Whereas the action of $L$ on $U$ is given by left multiplication the action of $L$ on $C(V)$ considered earlier in this section was given by Lie multiplication.

We consider the weights of the $L$-module $U$. In order to do this we identify the diagonal Cartan subalgebra $H$ of $L$.

Lemma 13.18 Under the above isomorphism $L \cong B_{l}$ the element $\left[v_{i} v_{-i}\right] \in L$ corresponds to the diagonal matrix

$$
i\left(\begin{array}{ccccccccc}
0 & & & & & & & & \\
\\
& \ddots & & & & & & & \\
\\
& & 0 & & & & & & \\
\\
& & & 4 & & & & & \\
\\
& & & & 0 & & & & \\
\\
& & & & & \ddots & & & \\
& & & & & & 0 & & \\
& & & & & & -4 & & \\
& & & & & & & & 0 \\
& & & & & & & & \\
& & & & & & & & \ddots
\end{array}\right)
$$

Proof. The matrix representation of $L$ comes from the $L$-module $V$ with basis $v_{0}, v_{1}, \ldots, v_{l}, v_{-1}, \ldots, v_{-l}$. Now

$$
\begin{aligned}
& {\left[\left[v_{i} v_{-i}\right], v_{0}\right]=0} \\
& {\left[\left[v_{i} v_{-i}\right], v_{j}\right]=4\left(v_{-i}, v_{j}\right) v_{i}-4\left(v_{i}, v_{j}\right) v_{-i}=\delta_{i j} \cdot 4 v_{i}} \\
& {\left[\left[v_{i} v_{-i}\right], v_{-j}\right]=4\left(v_{-i}, v_{-j}\right) v_{i}-4\left(v_{i}, v_{-j}\right) v_{-i}=-\delta_{i j} \cdot 4 v_{-i}}
\end{aligned}
$$

by Lemma 13.13 (i).

It follows from this lemma that the element

$$
h=\sum_{i=1}^{l} \frac{\lambda_{i}}{4\left[v_{i} v_{-i}\right]}
$$

is represented by the diagonal matrix

$$
\left(\begin{array}{lllllll}
0 & & & & & & \\
& \lambda_{1} & & & & & \\
& & \ddots & & & & \\
& & & \lambda_{l} & & & \\
& & & & -\lambda_{1} & & \\
& & & & & \ddots & \\
& & & & & & -\lambda_{l}
\end{array}\right)
$$

We recall from Section 8.3 that such matrices form a Cartan subalgebra $H$ of $L$.

We consider the action of $h$ on the $L$-module $U$. We have

$$
\begin{aligned}
{\left[v_{i} v_{-i}\right] } & =v_{i} v_{-i}-v_{-i} v_{i}=\frac{1}{2}\left(v_{i} v_{0} v_{0} v_{-i}-v_{-i} v_{0} v_{0} v_{i}\right) \\
& =-\frac{1}{2} u_{i} u_{-i}+\frac{1}{2} u_{-i} u_{i}
\end{aligned}
$$

Thus

$$
\begin{aligned}
{\left[v_{i} v_{-i}\right] u_{-j_{1}} \ldots u_{-j_{t}} u_{1} \ldots u_{l}=} & -\frac{1}{2} u_{i} u_{-i} \cdot u_{-j_{1}} \ldots u_{-j_{t}} u_{1} \ldots u_{l} \\
& +\frac{1}{2} u_{-i} u_{i} \cdot u_{-j_{1}} \ldots u_{-j_{t}} u_{1} \ldots u_{l}
\end{aligned} \quad \begin{aligned}
-2 u_{-j_{1}} \ldots u_{-j_{t}} u_{1} \ldots u_{l} \quad \text { if } i \in\left\{j_{1}, \ldots, j_{t}\right\} \\
2 u_{-j_{1}} \ldots u_{-j_{t}} u_{1} \ldots u_{l} \quad \text { if } i \notin\left\{j_{1}, \ldots, j_{t}\right\}
\end{aligned} .
$$

Thus

$$
h u_{-j_{1}} \ldots u_{-j_{t}} u_{1} \ldots u_{l}=\frac{1}{2}\left(\sum_{i=1}^{l} \varepsilon_{i} \lambda_{i}\right) u_{-j_{1}} \ldots u_{-j_{t}} u_{1} \ldots u_{l}
$$

where $\varepsilon_{i}=\left\{\begin{aligned}-1 & \text { if } i \in\left\{j_{1}, \ldots, j_{t}\right\} \\ 1 & \text { if } i \notin\left\{j_{1}, \ldots, j_{t}\right\}\end{aligned}\right.$. Let $\mu_{i} \in H^{*}$ be given by $\mu_{i}(h)=\lambda_{i}$. Then the weights of $L$ coming from the $L$-module $U$ are

$$
\frac{1}{2} \sum_{i=1}^{l} \varepsilon_{i} \mu_{i}
$$

for all possible choices of the signs $\varepsilon_{i}= \pm 1$. In particular the highest weight is $\frac{1}{2} \sum_{i=1}^{l} \mu_{i}$.

We recall from the proof of Theorem 13.9 that $\omega_{l}=\frac{1}{2}\left(\mu_{1}+\cdots+\mu_{l}\right)$. Thus $U$ has highest weight $\omega_{l}$. It follows that $U$ contains the spin module $L\left(\omega_{l}\right)$ as an irreducible direct summand. But

$$
\operatorname{dim} U=\operatorname{dim} L\left(\omega_{l}\right)=2^{l}
$$

Thus we have proved

Theorem 13.19 Let L be the simple Lie algebra of type $B_{l}$. Then the L-module $U$ constructed as above in the Clifford algebra is the spin module $L\left(\omega_{l}\right)$ of dimension $2^{l}$.

We now consider a second special case. This time let $V$ be a vector space with $\operatorname{dim} V=2 l$ and let $v_{1}, \ldots, v_{l}, v_{-1}, \ldots, v_{-l}$ be a basis of $V$. Consider the symmetric bilinear scalar product $V \times V \rightarrow \mathbb{C}$ determined by

$$
\left(v_{i}, v_{-i}\right)=1 \quad i=1, \ldots, l
$$

and all other scalar products of basis elements are 0 . The matrix of this scalar product is

$$
M=\left(\begin{array}{ll}
O & I_{l} \\
I_{l} & O
\end{array}\right)
$$

The condition of Lemma 13.16 implies that if $x \in L$ is represented by the matrix $X$ with respect to this basis of $V$ then

$$
X^{\mathrm{t}} M+M X=O
$$

Now $L$ has dimension $l(2 l-1)$ and acts faithfully on $V$. The Lie algebra of all $2 l \times 2 l$ matrices $X$ satisfying $X^{\mathrm{t}} M+M X=O$ is the simple Lie algebra $D_{l}$, by Section 8.2. Since $\operatorname{dim} D_{l}=l(2 l-1)$ we have $L=D_{l}$.

We again aim to find the two spin modules for $D_{l}$ inside the Clifford algebra $C(V)$. Let $U$ be the subspace of $C(V)$ spanned by all elements of form

$$
v_{-j_{1}} v_{-j_{2}} \ldots v_{-j_{t}} v_{1} \ldots v_{l}
$$

for all subsets $j_{1}<\cdots<j_{t}$ of $\{1, \ldots, l\}$. These elements are linearly independent, so form a basis for $U$. We have $\operatorname{dim} U=2^{l}$.

Lemma 13.20 $C(V) U \subset U$. Thus $U$ is a left ideal of $C(V)$.

Proof. $C(V)$ is generated by elements $v_{i}, v_{-i}$ and we have

$$
\begin{aligned}
& v_{i} \cdot v_{-j_{1}} \ldots v_{-j_{t}} v_{1} \ldots v_{l}= \begin{cases} \pm 2 v_{-j_{1}} \ldots \hat{v}_{-i} \ldots v_{-j_{t}} v_{1} \ldots v_{l} & \text { if } i \in\left\{j_{1}, \ldots, j_{t}\right\} \\
0 & \text { if } i \notin\left\{j_{1}, \ldots, j_{t}\right\}\end{cases} \\
& v_{-i} \cdot v_{-j_{1}} \ldots v_{-j_{t}} v_{1} \ldots v_{l}= \begin{cases}0 & \text { if } i \in\left\{j_{1}, \ldots, j_{t}\right\} \\
\pm v_{-j_{1}} \ldots v_{-i} \ldots v_{-j_{t}} v_{1} \ldots v_{l} & \text { if } i \notin\left\{j_{1}, \ldots, j_{t}\right\}\end{cases}
\end{aligned}
$$

Thus $v_{i} U \subset U$ and $v_{-i} U \subset U$, so $C(V) U \subset U$.
Let $U^{+}=U \cap C(V)^{+}$and $U^{-}=U \cap C(V)^{-}$. Then we have

$$
\begin{aligned}
& C(V)^{+} U^{+} \subset U \cap C(V)^{+}=U^{+} \\
& C(V)^{+} U^{-} \subset U \cap C(V)^{-}=U^{-} .
\end{aligned}
$$

Since $L \subset C(V)^{+}$it follows that

$$
L U^{+} \subset U^{+}, \quad L U^{-} \subset U^{-}
$$

Thus $U^{+}$and $U^{-}$are $L$-modules under left multiplication, with

$$
\operatorname{dim} U^{+}=\operatorname{dim} U^{-}=2^{l-1} .
$$

We shall show that these are the two spin modules for $L$.
We consider the weights of the $L$-modules $U^{+}, U^{-}$by identifying the diagonal Cartan subalgebra $H$ of $L$.

Lemma 13.21 Under the above isomorphism $L \cong D_{l}$ the element $\left[v_{i} v_{-i}\right] \in L$ corresponds to the diagonal matrix

$$
{ }_{-i}{ }^{-}\left(\begin{array}{ccccccccccc}
0 & & & & & & & & & & \\
& \ddots & & & & & & & & & \\
& & 0 & & & & & & & & \\
& & & 4 & & & & & & & \\
& & & & 0 & & & & & & \\
& & & & & \ddots & & & & & \\
& & & & & & 0 & & & & \\
& & & & & & & -4 & & & \\
& & & & & & & & 0 & & \\
& & & & & & & & & \ddots & \\
& & & & & & & & & & 0
\end{array}\right)
$$

Proof. The proof is the same as that for Lemma 13.18 with the first row and column omitted.

Thus the element

$$
h=\sum_{i=1}^{l} \frac{\lambda_{i}}{4\left[v_{i} v_{-i}\right]}
$$

is represented by the diagonal matrix

$$
\left(\begin{array}{cccccc}
\lambda_{1} & & & & & \\
& \ddots & & & & \\
& & \lambda_{l} & & & \\
& & & -\lambda_{1} & & \\
& & & & \ddots & \\
& & & & & -\lambda_{l}
\end{array}\right)
$$

We consider the action of $h$ on the $L$-modules $U^{+}$and $U^{-}$. We have

$$
\begin{aligned}
& {\left[v_{i} v_{-i}\right] v_{-j_{1}} \ldots v_{-j_{t}} v_{1} \ldots v_{l}} \\
& \quad=v_{i} v_{-i} v_{-j_{1}} \ldots v_{-j_{t}} v_{1} \ldots v_{l}-v_{-i} v_{i} v_{-j_{1}} \ldots v_{-j_{t}} v_{1} \ldots v_{l} \\
& \quad= \begin{cases}-2 v_{-j_{1}} \ldots v_{-j_{t}} v_{1} \ldots v_{l} & \text { if } i \in\left\{j_{1}, \ldots, j_{t}\right\} \\
2 v_{-j_{1}} \ldots v_{-j_{t}} v_{1} \ldots v_{l} & \text { if } i \notin\left\{j_{1}, \ldots, j_{t}\right\}\end{cases}
\end{aligned}
$$

since $v_{-i} v_{i}+v_{i} v_{-i}=21$. Thus

$$
h v_{-j_{1}} \ldots v_{-j_{t}} v_{1} \ldots v_{l}=\frac{1}{2}\left(\sum_{i=1}^{l} \varepsilon_{i} \lambda_{i}\right) v_{-j_{1}} \ldots v_{-j_{t}} v_{1} \ldots v_{l}
$$

where $\varepsilon_{i}=\left\{\begin{array}{ll}-1 & \text { if } i \in\left\{j_{1}, \ldots, j_{t}\right\} \\ 1 & \text { if } i \notin\left\{j_{1}, \ldots, j_{t}\right\} .\end{array}\right.$ As before let $\mu_{i} \in H^{*}$ be defined by

$$
\mu_{i}\left(\begin{array}{cccccc}
\lambda_{1} & & & & & \\
& \ddots & & & & \\
& & \lambda_{l} & & & \\
& & & -\lambda_{1} & & \\
& & & & \ddots & \\
& & & & & -\lambda_{l}
\end{array}\right)=\lambda_{i}
$$

Then the weights of the $L$-module $U$ are

$$
\frac{1}{2} \sum_{i=1}^{l} \varepsilon_{i} \mu_{i}
$$

for all possible choices of the signs $\varepsilon_{i}= \pm 1$.
If $l$ is even, the basis elements with $t$ even lie in $U^{+}$and those with $t$ odd in $U^{-}$. Thus the weights of $U^{+}$have an even number of $\varepsilon_{i}$ negative and those of $U^{-}$have an odd number negative. Since $\mu_{1} \succ \mu_{2} \succ \cdots \succ \mu_{l}$ the highest weight of $U^{+}$is $\frac{1}{2} \sum_{i=1}^{l} \mu_{i}$ and that of $U^{-}$is $\frac{1}{2}\left(\sum_{i=1}^{l-1} \mu_{i}-\mu_{l}\right)$.

If $l$ is odd we have the reverse situation in which $\frac{1}{2} \sum_{i=1}^{l} \mu_{i}$ is the highest weight of $U^{-}$and $\frac{1}{2}\left(\sum_{i=1}^{l-1} \mu_{i}-\mu_{l}\right)$ is the highest weight of $U^{+}$.

Now by the proof of Theorem 13.11 we have

$$
\begin{aligned}
& \frac{1}{2}\left(\mu_{1}+\cdots+\mu_{l-1}+\mu_{l}\right)=\omega_{l} \\
& \frac{1}{2}\left(\mu_{1}+\cdots+\mu_{l-1}-\mu_{l}\right)=\omega_{l-1}
\end{aligned}
$$

Thus we have proved
Theorem 13.22 Let $L$ be the simple Lie algebra of type $D_{l}$. Then the $L$-modules $U^{+}, U^{-}$are the spin modules of dimension $2^{l-1}$. If $l$ is even we have $U^{+}=L\left(\omega_{l}\right), U^{-}=L\left(\omega_{l-1}\right)$. If $l$ is odd we have $U^{+}=L\left(\omega_{l-1}\right), U^{-}=$ $L\left(\omega_{l}\right)$.

### 13.6 Fundamental modules for $C_{l}$

The fundamental weights $\omega_{1}, \ldots, \omega_{l}$ for a simple Lie algebra of type $C_{l}$ will be numbered according to the labelling of the Dynkin diagram


As before we shall use Theorem 13.1 to calculate $\operatorname{dim} L\left(\omega_{j}\right)$. We knows from Section 8.4 that the roots of $C_{l}$ have the following form. Let

$$
h=\left(\begin{array}{cccccc}
\lambda_{1} & & & & & \\
& \ddots & & & & \\
& & \lambda_{l} & & & \\
& & & -\lambda_{1} & & \\
& & & & \ddots & \\
& & & & & -\lambda_{l}
\end{array}\right) .
$$

Then the fundamental roots are

$$
\begin{aligned}
& \alpha_{i}(h)=\lambda_{i}-\lambda_{i+1} \quad \text { for } 1 \leq i \leq l-1 \\
& \alpha_{l}(h)=2 \lambda_{l}
\end{aligned}
$$

The full set of positive roots is given by

$$
\begin{array}{ll}
h \rightarrow \lambda_{i}-\lambda_{j} & \text { for } i<j \\
h \rightarrow \lambda_{i}+\lambda_{j} & \text { for } i<j \\
h \rightarrow 2 \lambda_{i} &
\end{array}
$$

where $i, j \in\{1, \ldots, l\}$. These positive roots can be expressed as combinations of fundamental roots as follows:

$$
\begin{aligned}
& \alpha_{i}+\cdots+\alpha_{j-1} \quad 1 \leq i<j \leq l \\
& \alpha_{i}+\cdots+\alpha_{j-1}+2 \alpha_{j}+\cdots+2 \alpha_{l-1}+\alpha_{l} \quad 1 \leq i<j \leq l \\
& 2 \alpha_{i}+\cdots+2 \alpha_{l-1}+\alpha_{l} \quad 1 \leq i \leq l
\end{aligned}
$$

The first two families are short roots and the third family are long roots. The weights $w_{i}$ are given by

$$
w_{1}=\cdots=w_{l-1}=1 \quad w_{l}=2
$$

According to Theorem 13.1 we have

$$
\operatorname{dim} L\left(\omega_{j}\right)=\prod_{\alpha \in \Phi^{+}} d_{\alpha}
$$

where $\alpha=\sum k_{i} \alpha_{i}$ and

$$
d_{\alpha}=\frac{\sum_{i=1}^{l} k_{i} w_{i}+k_{j} w_{j}}{\sum_{i=1}^{l} k_{i} w_{i}}
$$

We have $d_{\alpha}=1$ if $\alpha$ does not involve $\alpha_{j}$.
We first suppose that $j \in\{1, \ldots, l-1\}$. Then the positive roots involving $j$ are:

$$
\begin{array}{lc}
\alpha_{i}+\cdots+\alpha_{j}+\cdots+\alpha_{k-1} \quad 1 \leq i \leq j<k \leq l & \\
\alpha_{i}+\cdots+\alpha_{j}+\cdots+\alpha_{k-1}+2 \alpha_{k}+\cdots+2 \alpha_{l-1}+\alpha_{l} & 1 \leq i \leq j<k \leq l \\
\alpha_{i}+\cdots+\alpha_{k-1}+2 \alpha_{k}+\cdots+2 \alpha_{j}+\cdots+2 \alpha_{l-1}+\alpha_{l} & 1 \leq i<k \leq j \\
2 \alpha_{i}+\cdots+2 \alpha_{j}+\cdots+2 \alpha_{l-1}+\alpha_{l} \quad 1 \leq i \leq j &
\end{array}
$$

The values of $d_{\alpha}$ in these four cases are

$$
\frac{k-i+1}{k-i}, \quad \frac{2 l-i-k+3}{2 l-i-k+2}, \quad \frac{2 l-i-k+4}{2 l-i-k+2}, \quad \frac{l-i+2}{l-i+1}
$$

respectively. The product of all possible $d_{\alpha}$ in these four cases is

$$
\begin{aligned}
& \frac{(j+1)(j+2) \cdots l}{1 \cdot 2 \cdots \cdots l-j}, \quad \frac{(2 l-j+1)(2 l-j) \cdots(l+2)}{(2 l-2 j+1)(2 l-2 j) \cdots(l-j+2)}, \\
& \frac{(2 l+1) 2 l(2 l-1) \cdots(2 l-j+2)}{(2 l-j+2)(2 l-j+1) \cdots(2 l-2 j+3)}, \quad \frac{l+1}{l-j+1}
\end{aligned}
$$

respectively. The total product $\prod_{\alpha \in \Phi^{+}} d_{\alpha}$ is

$$
\frac{(2 l)!}{(2 l-j+2)!j!}(2 l+1)(2 l-2 j+2)
$$

This expression may be written in a more suggestive form by using the identity

$$
\binom{2 l}{j}-\binom{2 l}{j-2}=\frac{(2 l)!}{(2 l-j+2)!j!}(2 l+1)(2 l-2 j+2)
$$

Thus $\operatorname{dim} L\left(\omega_{j}\right)=\binom{2 l}{j}-\binom{2 l}{j-2}$ for $1 \leq j \leq l-1$.
We now suppose that $j=l$. The positive roots involving $\alpha_{l}$ are

$$
\begin{aligned}
& 2 \alpha_{i}+\cdots+2 \alpha_{l-1}+\alpha_{l} \quad 1 \leq i \leq l \\
& \alpha_{i}+\cdots+\alpha_{j-1}+2 \alpha_{j}+\cdots+2 \alpha_{l-1}+\alpha_{l} \quad 1 \leq i<j \leq l
\end{aligned}
$$

The first family are long roots and the second short roots. The values of $1 d_{\alpha}$ in these two cases are

$$
\frac{l-i+2}{l-i+1}, \quad \frac{2 l-i-j+4}{2 l-i-j+2}
$$

respectively. The product of all possible $d_{\alpha}$ in these cases is

$$
l+1, \quad \frac{(2 l+1)(2 l) \cdots(l+2)}{(l+2)(l+1) \cdots 3}
$$

respectively, and the total product $\prod_{\alpha \in \Phi^{+}} d_{\alpha}$ is $\frac{(2 l+1)!2}{(l+2)!!!}$.

By using the identity

$$
\binom{2 l}{l}-\binom{2 l}{l-2}=\frac{(2 l+1)!2}{(l+2)!l!}
$$

we see that

$$
\operatorname{dim} L\left(\omega_{l}\right)=\binom{2 l}{l}-\binom{2 l}{l-2}
$$

Thus we have shown

Proposition 13.23 The dimensions of the fundamental modules for the simple Lie algebra of type $C_{l}$ are


### 13.7 Contraction maps

We shall now identify the fundamental modules whose dimensions we have obtained. We begin with $L\left(\omega_{1}\right)$.

Proposition 13.24 The natural $2 l$-dimensional $C_{l}$-module is isomorplic to $L\left(\omega_{1}\right)$.

Proof. Let $V$ be the natural $C_{l}$-module. Let

$$
h=\left(\begin{array}{cccccc}
\lambda_{1} & & & & & \\
& \ddots & & & & \\
& & \lambda_{l} & & & \\
& & & -\lambda_{1} & & \\
& & & & \ddots & \\
& & & & & -\lambda_{l}
\end{array}\right)
$$

Then the weights of $V$ are $\mu_{1}, \ldots, \mu_{l},-\mu_{1}, \ldots,-\mu_{l}$ where $\mu_{i}(h)=\lambda_{i}$. Since

$$
\begin{aligned}
& \mu_{i}-\mu_{i+1}=\alpha_{i} \quad 1 \leq i \leq l-1 \\
& 2 \mu_{l}=\alpha_{l}
\end{aligned}
$$

we have

$$
\mu_{1} \succ \mu_{2} \succ \cdots \succ \mu_{l} \succ 0 .
$$

Thus the highest weight of $V$ is $\mu_{1}$. We have

$$
\begin{aligned}
\mu_{1} & =\alpha_{1}+\cdots+\alpha_{l-1}+\frac{1}{2} \alpha_{l} \\
\mu_{2} & =\alpha_{2}+\cdots+\alpha_{l-1}+\frac{1}{2} \alpha_{l} \\
& \vdots \\
\mu_{l} & =\frac{1}{2} \alpha_{l} .
\end{aligned}
$$

We also have $\alpha_{i}=\sum_{j} A_{j i} \omega_{j}$ which in type $C_{l}$ gives

$$
\begin{aligned}
& \alpha_{1}=2 \omega_{1}-\omega_{2} \\
& \alpha_{2}=-\omega_{1}+2 \omega_{2}-\omega_{3} \\
& \vdots \\
& \alpha_{l-1}=-\omega_{l-2}+2 \omega_{l-1}-\omega_{l} \\
& \alpha_{l}=-2 \omega_{l-1}+2 \omega_{l} .
\end{aligned}
$$

It follows that

$$
\begin{aligned}
& \mu_{1}=\omega_{1} \\
& \mu_{2}=-\omega_{1}+\omega_{2} \\
& \vdots \\
& \mu_{l}=-\omega_{l-1}+\omega_{l} .
\end{aligned}
$$

Thus $V$ is a $C_{l}$-module with highest weight $\omega_{1}$. It therefore contains $L\left(\omega_{1}\right)$ as an irreducible component. However,

$$
\operatorname{dim} V=\operatorname{dim} L\left(\omega_{1}\right)=2 l
$$

thus $V$ is irreducible and isomorphic to $L\left(\omega_{1}\right)$.
We now consider the fundamental modules $L\left(\omega_{j}\right)$ for $j \geq 2$. We have

$$
\operatorname{dim} L\left(\omega_{j}\right)=\binom{2 l}{j}-\binom{2 l}{j-2}
$$

This suggests that we should look for $L\left(\omega_{j}\right)$ as a submodule of the exterior power $\Lambda^{j} V$. The key idea is to find a homomorphism of $C_{l}$-modules from $\Lambda^{j} V$ into $\Lambda^{j-2} V$, called a contraction map.

Proposition 13.25 Let $v, v^{\prime} \rightarrow\left(v, v^{\prime}\right)$ be the skew-symmetric bilinear map $V \times V \rightarrow \mathbb{C}$ given by the matrix

$$
M=\left(\begin{array}{cc}
O & I_{l} \\
-I_{l} & O
\end{array}\right)
$$

Then there is a unique homomorphism of $C_{l}$-modules

$$
\theta: \Lambda^{j} V \rightarrow \Lambda^{j-2} V
$$

satisfying the condition

$$
\begin{align*}
\theta\left(u_{1} \wedge \cdots \wedge u_{j}\right)= & \sum_{r<s}(-1)^{r+s-1}\left(u_{r}, u_{s}\right) u_{1} \wedge \cdots \wedge \hat{u}_{r} \wedge \cdots \wedge \\
& \hat{u}_{s} \wedge \cdots \wedge u_{j} \text { forall } u_{1}, \ldots, u_{j} \in V
\end{align*}
$$

Here as usual the notation $\hat{u}_{r}, \hat{u}_{s}$ means that those terms are omitted.
Proof. It is clear that if such a map $\theta$ exists it will be unique. To prove the existence let $v_{1}, \ldots, v_{2 l}$ be a basis of $V$. Then there is a unique linear map $\theta$ satisfying

$$
\theta\left(v_{i_{1}} \wedge \cdots \wedge v_{i_{j}}\right)=\sum_{r<s}(-1)^{r+s-1}\left(v_{i_{r}}, v_{i_{s}}\right) v_{i_{1}} \wedge \cdots \wedge \hat{v}_{i_{r}} \wedge \cdots \wedge \hat{v}_{i_{s}} \wedge \cdots \wedge v_{i_{j}}
$$

for all $i_{1}, \ldots, i_{j} \in\{1, \ldots, 2 l\}$ with $i_{1}<\cdots<i_{j}$. We show this map has the required properties. Since both sides of equation $(\dagger)$ are linear in $u_{1}, \ldots, u_{j}$ it will be sufficient to prove it when each $u_{k}$ is one of the basis elements of $V$. If the same basis element appears twice both sides of $(\dagger)$ are 0 . Thus we may assume the basis elements are all distinct. They may not occur in increasing order, thus we must show that the above formula defining $\theta$ remains valid if the factors $v_{i_{1}}, \ldots, v_{i_{j}}$ are permuted. In fact it is sufficient to see this if we transpose two consecutive terms $v_{i_{k}}, v_{i_{k+1}}$. When we carry out such a transposition the expression $v_{i_{1}} \wedge \cdots \wedge v_{i_{j}}$ changes in sign. We show that each term

$$
(-1)^{r+s-1}\left(v_{i_{r}}, v_{i_{s}}\right) v_{i_{1}} \wedge \cdots \wedge \hat{v}_{i_{r}} \wedge \cdots \wedge \hat{v}_{i_{s}} \wedge \cdots \wedge v_{i_{j}}
$$

changes in sign also. If neither of $r, s$ lie in $\{k, k+1\}$ the term

$$
v_{i_{1}} \wedge \cdots \wedge \hat{v}_{i_{r}} \wedge \cdots \wedge \hat{v}_{i_{s}} \wedge \cdots \wedge \hat{v}_{i_{j}}
$$

will change in sign when we make the transposition. If just one of $r, s$ lies in $\{k, k+1\}$ the term $(-1)^{r+s-1}$ will change in sign when the transposition is made. Finally if $r=k, s=k+1$ the term $\left(v_{i_{r}}, v_{i_{s}}\right)$ changes in sign, since the bilinear map is skew-symmetric. This shows that the linear map $\theta$ we have defined satisfies $(\dagger)$.

It remains to show that $\theta$ is a homomorphism of $C_{l}$-modules. Let $x$ lie in the Lie algebra $C_{l}$. Then

$$
\begin{gathered}
x \theta\left(u_{1} \wedge \cdots \wedge u_{j}\right)=x \sum_{r<s}(-1)^{r+s-1}\left(u_{r}, u_{s}\right) u_{1} \wedge \cdots \wedge \hat{u}_{r} \wedge \cdots \wedge \hat{u}_{s} \wedge \cdots \wedge u_{j} \\
=\sum_{r<s} \sum_{\substack{k \\
k \neq r, k \neq s}}(-1)^{r+s-1}\left(u_{r}, u_{s}\right) u_{1} \wedge \cdots \wedge x u_{k} \wedge \cdots \wedge \hat{u}_{r} \wedge \cdots \\
\wedge \hat{u}_{s} \wedge \cdots \wedge u_{j}
\end{gathered}
$$

On the other hand we have

$$
\begin{aligned}
\theta x\left(u_{1} \wedge\right. & \left.\cdots \wedge u_{j}\right) \\
= & \theta \sum_{k} u_{1} \wedge \cdots \wedge x u_{k} \wedge \cdots \wedge u_{j}=x \theta\left(u_{1} \wedge \cdots \wedge u_{j}\right) \\
& +\sum_{\substack{k \\
k<s}} \sum_{s}(-1)^{k+s-1}\left(x u_{k}, u_{s}\right) u_{1} \wedge \cdots \wedge \hat{u}_{k} \wedge \cdots \wedge \hat{u}_{s} \wedge \cdots \wedge u_{j} \\
& +\sum_{k} \sum_{r}(-1)^{r+k-1}\left(u_{r}, x u_{k}\right) u_{1} \wedge \cdots \wedge \hat{u}_{r} \wedge \cdots \wedge \hat{u}_{k} \wedge \cdots \wedge u_{j}
\end{aligned}
$$

Renaming the suffixes we see that the last two sums cancel since

$$
\left(x u_{k}, u_{s}\right)+\left(u_{k}, x u_{s}\right)=0
$$

This condition is equivalent to

$$
X^{\mathrm{t}} M+M X=O
$$

where $X$ is the matrix representing $x$ on $V$, and we recall from Section 8.4 that the simple Lie algebra $C_{l}$ satisfies this condition. It follows that

$$
\theta x\left(u_{1} \wedge \cdots \wedge u_{j}\right)=x \theta\left(u_{1} \wedge \cdots \wedge u_{j}\right)
$$

and so $\theta$ is a homomorphism of $C_{l}$-modules.
This homomorphism $\theta: \Lambda^{j} V \rightarrow \Lambda^{j-2} V$ will be called a contraction map.
Now the weights of $\Lambda^{j} V$ are sums of $j$ distinct weights of $V$. By the proof of Proposition 13.24 the weights of $V$ are

$$
\begin{aligned}
\omega_{1} & \succ-\omega_{1}+\omega_{2} \succ-\omega_{2}+\omega_{3} \succ \cdots \succ-\omega_{l-1}+\omega_{l} \\
& \succ \omega_{l-1}-\omega_{l} \succ \cdots \succ \omega_{2}-\omega_{3} \succ \omega_{1}-\omega_{2} \succ-\omega_{1} .
\end{aligned}
$$

Thus if $j \leq l$ the highest weight of $\Lambda^{j} V$ is $\omega_{j}$. Similarly the highest weight of $\Lambda^{j-2} V$ is $\omega_{j-2}$. Since $\omega_{j}>\omega_{j-2}$ we see that $\omega_{j}$ is not a weight of $\Lambda^{j-2} V$. Since $\omega_{j}$ is the highest weight of $\Lambda^{j} V$ the module $L\left(\omega_{j}\right)$ must be an irreducible
direct summand of $\Lambda^{j} V$. On the other hand $L\left(\omega_{j}\right)$ cannot be a submodule of $\Lambda^{j-2} V$, as $\omega_{j}$ is not a weight of this module. Thus $L\left(\omega_{j}\right)$ must lie in the kernel of the contraction map $\theta$.

We shall show subsequently that when $j \leq l$ the contraction map $\theta: \Lambda^{j} V \rightarrow$ $\Lambda^{j-2} V$ is surjective. It will follow that

$$
\operatorname{dim}(\operatorname{ker} \theta)=\binom{2 l}{j}-\binom{2 l}{j-2}=\operatorname{dim} L\left(\omega_{j}\right)
$$

and therefore that $L\left(\omega_{j}\right)=\operatorname{ker} \theta$. This will identify the irreducible module $L\left(\omega_{j}\right)$ as the submodule of $\Lambda^{j} V$ which is the kernel of the contraction map $\theta$.

Let $v_{1}, \ldots, v_{l}, v_{-1}, \ldots, v_{-l}$ be the natural basis of $V$ with respect to which the skew-symmetric bilinear form is given by

$$
\begin{aligned}
& \left(v_{i}, v_{-i}\right)=1 \quad 1 \leq i \leq l \\
& \left(v_{-i}, v_{i}\right)=-1
\end{aligned}
$$

and all other scalar products zero. Let $W$ be the subspace of $V$ spanned by $v_{1}, \ldots, v_{l}$ and $W^{-}$the subspace spanned by $v_{-1}, \ldots, v_{-l}$. Then $W, W^{-}$are isotropic subspaces of $V$, i.e. the skew-symmetric form restricted to $W$ and $W^{-}$is identically zero. Also we have $V=W \oplus W^{-}$. It follows that

$$
\Lambda^{j} V=\bigoplus_{a+b=j}\left(\Lambda^{a} W \otimes \Lambda^{b} W^{-}\right)
$$

The contraction map $\theta: \Lambda^{j} V \rightarrow \Lambda^{j-2} V$ satisfies

$$
\theta\left(\Lambda^{a} W \otimes \Lambda^{b} W^{-}\right) \subset \Lambda^{a-1} W \otimes \Lambda^{b-1} W^{-}
$$

since a basis element in $W$ has a non-zero scalar product only with a basis element in $W^{-}$. Thus in order to show that $\theta: \Lambda^{j} V \rightarrow \Lambda^{j-2} V$ is surjective for $j \leq l$ it will be sufficient to show that

$$
\theta: \Lambda^{a} W \otimes \Lambda^{b} W^{-} \rightarrow \Lambda^{a-1} W \otimes \Lambda^{b-1} W^{-}
$$

is surjective whenever $a+b \leq l$.
For each subset $I \subset\{1, \ldots, l\}$ we define $v_{I}=v_{i_{1}} \wedge \cdots \wedge v_{i_{k}}$ where $I=$ $\left\{i_{1}, \ldots, i_{k}\right\}$ with $i_{1}<\cdots<i_{k}$. We also define $v_{-I}=v_{-i_{1}} \wedge \cdots \wedge v_{-i_{k}}$. Then any basis element of $\Lambda^{a-1} W \otimes \Lambda^{b-1} W^{-}$can be written in the form

$$
\pm\left(v_{X} \wedge v_{T}\right) \otimes\left(v_{-T} \wedge v_{-Y}\right)
$$

for some subsets $T, X, Y$ of $\{1, \ldots, l\}$ with
$T \cap X=\phi, \quad T \cap Y=\phi, \quad X \cap Y=\phi, \quad|X|+|T|=a-1, \quad|Y|+|T|=b-1$.

We write $|T|=r$. Since $a+b \leq l$ we have $|X|+|Y|+2 r+2 \leq l$, that is $l-|X|-$ $|Y| \geq 2 r+2$. Thus it is possible to choose a subset $S$ of $\{1, \ldots, l\}$ such that

$$
|S|=2 r+1, \quad S \cap X=\phi, \quad S \cap Y=\phi, \quad S \supset T .
$$

We can now describe an element of $\Lambda^{a} W \otimes \Lambda^{b} W^{-}$which maps under $\theta$ to a non-zero multiple of $\left(v_{X} \wedge v_{T}\right) \otimes\left(v_{-T} \wedge v_{-Y}\right)$.

Proposition 13.26 Suppose subsets $T, X, Y, S$ of $\{1, \ldots, l\}$ are chosen as above, and let $\theta: \Lambda^{j} V \rightarrow \Lambda^{j-2} V$ be the contraction map. Then

$$
\begin{aligned}
& \theta\binom{\sum_{i=0}^{r}(-1)^{i} i!(r-i)!\sum_{\substack{U \\
U \backslash S \\
|U| r+1 \\
|U \cap T|=i}}\left(v_{X} \wedge v_{U}\right) \otimes\left(v_{-U} \wedge v_{-Y}\right)}{\quad=(r+1)!\left(v_{X} \wedge v_{T}\right) \otimes\left(v_{-T} \wedge v_{-Y}\right) .}
\end{aligned}
$$

Consequently the map

$$
\theta: \Lambda^{a} W \otimes \Lambda^{b} W^{-} \rightarrow \Lambda^{a-1} W \otimes \Lambda^{b-1} W^{-}
$$

is surjective when $a+b \leq l$.
Proof. We note that $S \supset T,|S|=2 r+1,|T|=r$ and that we are summing over all subsets $U$ of $S$ with $|U|=r+1$ and $|U \cap T|=i$. Since $|X|+|T|=a-1$ and $|Y|+|T|=b-1$ we have $|X|+|U|=a$ and $|Y|+|U|=b$. Thus the left-hand side lies in $\Lambda^{a} W \otimes \Lambda^{b} W^{-}$.
By definition of $\theta$ we have

$$
\theta\left(v_{u} \otimes v_{-u}\right)=(-1)^{r} \sum_{\substack{R \\ R U \\|R|=r}} v_{R} \otimes v_{-R}
$$

where the right-hand side involves a sum over all $r$-element subsets $R$ of $U$. Thus

$$
\begin{aligned}
\theta\left(\sum_{\substack{U \\
|U \cap T|=i}} v_{U} \otimes v_{-U}\right) & =(-1)^{r} \sum_{\substack{U \\
|U|=r+1 \\
|U \cap T|=i}}\left(\sum_{\substack{R \\
|R|=r \\
R \subset U}} v_{R} \otimes v_{-R}\right) \\
& =(-1)^{r} \sum_{\substack{R \\
|R|=r}}\left(\begin{array}{c}
\substack{U \\
R \subset U \\
|U|=r+1 \\
|U \cap T|=i}
\end{array}\right)
\end{aligned}
$$

Since $|U \cap T|=i$ and $R$ is obtained from $U$ by omitting one element we have $|R \cap T|=i$ or $|R \cap T|=i-1$. We split the sum according to those two possibilities. Thus

$$
\begin{aligned}
& \theta\left(\sum_{\substack{U \\
|U \cap T|=i}} v_{U} \otimes v_{-U}\right) \\
& \left.\quad=(-1)^{r} \sum_{\substack{R \\
|R|=r \\
|R \cap T|=i \\
\hline}}\left(\sum_{\substack{U \\
R \subset U \\
|U|=r+1 \\
|U \cap T|=i}} 1\right) v_{R} \otimes v_{-R}+(-1)^{r} \sum_{\substack{R \\
|R|=r \\
|R \cap T|=i-1}} \sum_{\substack{U \\
|R \subset U\\
| U|=+1\\
| U \cap T \mid=i}} 1\right) v_{R} \otimes v_{-R} \\
& =(-1)^{r} \sum_{\substack{R \\
|R|=r \\
|R \cap T|=i}}(i+1) v_{R} \otimes v_{-R}+(-1)^{r} \sum_{\substack{R \\
|R|=r \\
|R \cap T|=i-1}}(r+1-i) v_{R} \otimes v_{-R}
\end{aligned}
$$

since in the first case the additional element of $U$ can be chosen in $i+1$ ways and in the second case in $r+1-i$ ways. Thus

$$
\begin{aligned}
& \theta\left(\sum_{i=0}^{r}(-1)^{i} i!(r-i)!\sum_{\substack{U \\
|U \cap T|=i}} v_{U} \otimes v_{-U}\right) \\
& =(-1)^{r} \sum_{i=0}^{r}(-1)^{i} i!(r-i)!\sum_{\substack{R \\
|R|=r \\
|R \cap T|=i}}(i+1) v_{R} \otimes v_{-R} \\
& +(-1)^{r} \sum_{i=0}^{r}(-1)^{i} i!(r-i)!\sum_{\substack{R \\
|R|=r \\
|R \cap T|=i-1}}(r+1-i) v_{R} \otimes v_{-R} .
\end{aligned}
$$

We rename the variable $i$ in the second sum to give

$$
\begin{aligned}
& (-1)^{r} \sum_{i=0}^{r}(-1)^{i} i!(r-i)!\sum_{\substack{R \\
|R|=r \\
|R \cap T|=i}}(i+1) v_{R} \otimes v_{-R} \\
& +(-1)^{r} \sum_{i=-1}^{r-1}(-1)^{i+1}(i+1)!(r-i-1)!\sum_{\substack{R \\
|R|=r \\
|R \cap T|=i}}(r-i) v_{R} \otimes v_{-R}
\end{aligned}
$$

$$
\begin{aligned}
= & (-1)^{r} \sum_{i=0}^{r-1}(-1)^{i}(i+1)!(r-i)!\sum_{\substack{R \\
|R|=r \\
|R \cap T|=i}}(1-1) v_{R} \otimes v_{-R} \\
& +(r+1)!\sum_{\substack{R \\
|R=r\\
| R \cap T \mid=r}} v_{R} \otimes v_{-R} \\
= & (r+1)!v_{T} \otimes v_{-T} .
\end{aligned}
$$

We now consider

$$
\theta\left(\sum_{i=0}^{r}(-1)^{i} i!(r-i)!\sum_{\substack{U \\ U \subset S \\|U|=r+1 \\|U \cap T|=i}}\left(v_{X} \wedge v_{U}\right) \otimes\left(v_{-U} \wedge v_{-Y}\right)\right)
$$

Since the $v_{i}$ for $i \in X$ and the $v_{-i}$ for $i \in Y$ have scalar product 0 with all factors in the above product they are not involved in any contraction. Thus

$$
\begin{aligned}
& \theta\left(\sum_{i=0}^{r}(-1)^{i} i!(r-i)!\sum_{|U \cap T|=i}\left(v_{X} \wedge v_{U}\right) \otimes\left(v_{-U} \wedge v_{-Y}\right)\right) \\
& \quad=v_{X} \wedge \theta\left(\sum_{i=0}^{r}(-1)^{i} i!(r-i)!\sum_{|U \cap T|=i} v_{U} \otimes v_{-U}\right) \wedge v_{-Y} \\
& \quad=v_{X} \wedge\left((r+1)!v_{T} \otimes v_{-T}\right) \wedge v_{-Y} \\
& \quad=(r+1)!\left(v_{X} \wedge v_{T}\right) \otimes\left(v_{-T} \wedge v_{-Y}\right)
\end{aligned}
$$

Corollary 13.27 The contraction map $\theta: \Lambda^{j} V \rightarrow \Lambda^{j-2} V$ is surjective when $j \leq l$.

The surjectivity of $\theta$ enables us to identify the fundamental modules $L\left(\omega_{j}\right)$.
Theorem 13.28 The fundamental modules $L\left(\omega_{j}\right)$ for the simple Lie algebra $C_{l}$ are given as follows.
(a) $L\left(\omega_{1}\right)$ is the natural 2l-dimensional $C_{l}$-module $V$.
(b) For $2 \leq j \leq l, L\left(\omega_{j}\right)$ is the submodule of $\Lambda^{j} V$ given by the kernel of the contraction map $\theta: \Lambda^{j} V \rightarrow \Lambda^{j-2} V$.

Proof. (a) was shown in Proposition 13.24. We also pointed out earlier that $L\left(\omega_{j}\right)$ is a submodule of $\Lambda^{j} V$ contained in the kernel of $\theta$. Since $\theta: \Lambda^{j} V \rightarrow$ $\Lambda^{j-2} V$ is surjective for $j \leq l$ we have

$$
\operatorname{dim} \operatorname{ker} \theta=\binom{2 l}{j}-\binom{2 l}{j-2}
$$

and this is equal to $\operatorname{dim} L\left(\omega_{j}\right)$ by Proposition 13.23. It follows that $L\left(\omega_{j}\right)=$ $\operatorname{ker} \theta$.

### 13.8 Fundamental modules for exceptional algebras

By applying Theorem 13.1 to the exceptional simple Lie algebras and making use of the information about their root systems available in Sections 8.5, 8.6 and 8.7 we can show that the dimensions of the fundamental modules for these algebras are as shown. We omit the details.


We shall show in each case how to obtain the fundamental module of smallest dimension. We begin by obtaining a 27-dimensional fundamental module for $E_{6}$.

Proposition 13.29 (a) The number of positive roots of $E_{7}$ not in $E_{6}$ is 27.
(b) The subspace $V$ of $E_{7}$ spanned by vectors $e_{\alpha}$ for such roots is a 27dimensional fundamental $E_{6}$-module.

Proof. We recall from Section 8.7 that the fundamental roots of $E_{7}$ are given by

and that the full set of roots of $E_{7}$ is

$$
\begin{aligned}
& \pm \beta_{i} \pm \beta_{j} \quad i \neq j \quad i, j \in\{2,3,4,5,6,7\} \\
& \pm\left(\beta_{1}+\beta_{8}\right) \\
& \frac{1}{2} \sum \varepsilon_{i} \beta_{i} \quad \varepsilon_{i} \in\{1,-1\}, \quad \prod \varepsilon_{i}=1, \quad \varepsilon_{1}=\varepsilon_{8}
\end{aligned}
$$

The positive roots are

$$
\begin{array}{ll}
\beta_{i}-\beta_{j} \quad i \neq j \quad i, j \in\{2,3,4,5,6,7\} \\
\beta_{i}+\beta_{j} & i \neq j \quad i, j \in\{2,3,4,5,6,7\} \\
-\beta_{1}-\beta_{8} & \\
\frac{1}{2} \sum \varepsilon_{i} \beta_{i} \quad \varepsilon_{i} \in\{1,-1\}, \quad \prod \varepsilon_{i}=1, \quad \varepsilon_{1}=\varepsilon_{8}=-1 .
\end{array}
$$

The positive roots of $E_{7}$ which are not roots of $E_{6}$ are

$$
\begin{array}{ll}
\beta_{2}-\beta_{j} & j \in\{3,4,5,6,7\} \\
\beta_{2}+\beta_{j} & j \in\{3,4,5,6,7\} \\
-\beta_{1}-\beta_{8} & \\
\frac{1}{2} \sum \varepsilon_{i} \beta_{i} & \prod \varepsilon_{i}=1, \quad \varepsilon_{1}=\varepsilon_{8}=-1, \quad \varepsilon_{2}=1
\end{array}
$$

The number of such roots is 27 .
Now let $V$ be the subspace of $E_{7}$ spanned by the root vectors $e_{\alpha}$ for such roots $\alpha$. Then $\operatorname{dim} V=27$.

Now $E_{7}$ may be regarded as an $E_{7}$-module giving the adjoint representation. In particular $E_{7}$ may be regarded as an $E_{6}$-module. We observe that $V$ is an $E_{6}$-submodule. To see this it is sufficient to show that $\left[e_{\alpha} e_{\beta}\right] \in V$ for all $\alpha \in \Phi\left(E_{6}\right), \beta \in \Phi^{+}\left(E_{7}\right)-\Phi^{+}\left(E_{6}\right)$. We have

$$
\left[e_{\alpha} e_{\beta}\right]= \begin{cases}N_{\alpha, \beta} e_{\alpha+\beta} & \text { if } \alpha+\beta \in \Phi\left(E_{7}\right) \\ 0 & \text { otherwise }\end{cases}
$$

Suppose $\alpha+\beta \in \Phi\left(E_{7}\right)$. Since $\beta$ is not a root of $E_{6}, \beta$ will involve the fundamental root of $E_{7}$ not in $E_{6}$, and since $\beta$ is positive this fundamental root will have positive coefficient in $\beta$. It will therefore have positive coefficient in $\alpha+\beta$, and so $\alpha+\beta \in \Phi^{+}\left(E_{7}\right)$. We claim that $\alpha+\beta \notin \Phi\left(E_{6}\right)$. Suppose to the contrary that $\alpha+\beta \in \Phi\left(E_{6}\right)$. Then $-\alpha \in \Phi\left(E_{6}\right)$ and

$$
\left[e_{\alpha+\beta} e_{-\alpha}\right]=N_{\alpha+\beta,-\alpha} e_{\beta}
$$

Since $N_{\alpha, \beta} \neq 0$ it follows from Proposition 7.1 that $N_{\alpha+\beta,-\alpha} \neq 0$ and so $\beta \in \Phi\left(E_{6}\right)$, a contradiction. Hence $\alpha+\beta \in \Phi^{+}\left(E_{7}\right)-\Phi^{+}\left(E_{6}\right)$ and $V$ is an $E_{6}$-module.

In order to determine the highest weight of $V$ it is convenient to use the linear function

$$
h: \sum_{i=1}^{8} \mathbb{R} \boldsymbol{\beta}_{i} \rightarrow \mathbb{R}
$$

determined by the property that $h\left(\alpha_{i}\right)=1$ for each fundamental root $\alpha_{i}$ of $E_{8}$. Thus

$$
\begin{aligned}
& h\left(\beta_{i}-\beta_{i+1}\right)=1 \quad \text { for } i \in\{1, \ldots, 6\} \\
& h\left(\beta_{6}+\beta_{7}\right)=1 \\
& h\left(-\frac{1}{2} \sum_{i=1}^{8} \beta_{i}\right)=1
\end{aligned}
$$

Hence we have

$$
\begin{aligned}
& h\left(\beta_{1}\right)=6, \quad h\left(\beta_{2}\right)=5, \quad h\left(\beta_{3}\right)=4, \quad h\left(\beta_{4}\right)=3, \quad h\left(\beta_{5}\right)=2 \\
& \quad h\left(\beta_{6}\right)=1, \quad h\left(\beta_{7}\right)=0, \quad h\left(\beta_{8}\right)=-23
\end{aligned}
$$

Of our 27 roots the one with the highest $h$-value is $-\beta_{1}-\beta_{8}$. This must therefore be a highest weight of $V$. Now the fundamental roots of $E_{6}$ are

and $-\beta_{1}-\beta_{8}$ is orthogonal to all of them except $-\frac{1}{2} \sum_{i=1}^{8} \beta_{i}$. Moreover the scalar product $\{$,$\} satisfies$

$$
\left\{-\beta_{1}-\beta_{8},-\frac{1}{2} \sum_{i=1}^{8} \beta_{i}\right\}=\frac{1}{2}\left\{-\frac{1}{2} \sum_{i=1}^{8} \beta_{i},-\frac{1}{2} \sum_{i=1}^{8} \beta_{i}\right\}
$$

thus $-\beta_{1}-\beta_{8}$ is the fundamental weight $\omega_{8}$. Hence $L\left(\omega_{8}\right)$ is an irreducible direct summand of $V$. Since

$$
\operatorname{dim} L\left(\omega_{8}\right)=\operatorname{dim} V=27
$$

we deduce that $V=L\left(\omega_{8}\right)$.
In order to obtain the other 27-dimensional fundamental $E_{6}$-module we introduce the dual module. We recall that, given any $L$-module $V$, the dual space $V^{*}$ of linear maps from $V$ to $\mathbb{C}$ may be made into an L-module by the rule

$$
(x f) v=-f(x v) \quad x \in L, \quad f \in V^{*}, \quad v \in V
$$

The weights of $V^{*}$ are the negatives of the weights of $V$. In the case of the 27-dimensional $E_{6}$-module $V$ above, the highest weight of $V^{*}$ is the negative of the lowest weight of $V$. The lowest weight of $V$ is the one with the smallest value of $h$, i.e. $\beta_{2}-\beta_{3}$. Thus the highest weight of $V^{*}$ is $\beta_{3}-\beta_{2}$. This is orthogonal to all fundamental roots of $E_{6}$ except for $\alpha_{3}=\beta_{3}-\beta_{4}$. Since

$$
\left\{\beta_{3}-\beta_{2}, \beta_{3}-\beta_{4}\right\}=\frac{1}{2}\left\{\beta_{3}-\beta_{4}, \beta_{3}-\beta_{4}\right\}
$$

we deduce that $\beta_{3}-\beta_{2}=\omega_{3}$. Hence $V^{*}=L\left(\omega_{3}\right)$.
Now the weight of $V$ with second highest value of $h$ is $\frac{1}{2}\left(-\beta_{1}+\beta_{2}+\beta_{3}+\right.$ $\left.\beta_{4}+\beta_{5}+\beta_{6}+\beta_{7}-\beta_{8}\right)$ and the third highest is $\frac{1}{2}\left(-\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}\right.$ $-\beta_{6}-\beta_{7}-\beta_{8}$ ). Thus the highest weight of $\Lambda^{2} V$ is

$$
\begin{gathered}
\left(-\beta_{1}-\beta_{8}\right)+\frac{1}{2}\left(-\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}+\beta_{7}-\beta_{8}\right) \\
=\frac{1}{2}\left(-3 \beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}+\beta_{7}-3 \beta_{8}\right)
\end{gathered}
$$

By considering the scalar products $\{$,$\} of this weight with the fundamental$ roots of $E_{6}$ we see that this weight is $\omega_{7}$. Since

$$
\operatorname{dim} L\left(\omega_{7}\right)=\binom{27}{2}=\operatorname{dim} \Lambda^{2} V
$$

we deduce $\Lambda^{2} V=L\left(\omega_{7}\right)$.
Similarly the highest weight of $\Lambda^{3} V$ is

$$
\begin{gathered}
\left(-\beta_{1}-\beta_{8}\right)+\frac{1}{2}\left(-\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}+\beta_{7}-\beta_{8}\right) \\
+\frac{1}{2}\left(-\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}-\beta_{6}-\beta_{7}-\beta_{8}\right) \\
=-2 \beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}-2 \beta_{8}
\end{gathered}
$$

We check by computing scalar products $\{$,$\} that this is the weight \omega_{5}$ of $E_{6}$. Since

$$
\operatorname{dim} L\left(\omega_{5}\right)=\binom{27}{3}=\operatorname{dim} \Lambda^{3} V
$$

we deduce $\Lambda^{3} V=L\left(\omega_{5}\right)$.
It may be shown similarly that

$$
\Lambda^{2} V^{*}=L\left(\omega_{4}\right) \quad \text { and } \quad \Lambda^{3} V^{*}=\Lambda^{3} V=L\left(\omega_{5}\right)
$$

Finally $L\left(\omega_{6}\right)$ is the adjoint module. Thus the fundamental $E_{6}$-modules are


We now consider the simple Lie algebra $E_{7}$ and obtain a 56-dimensional fundamental module. The idea is similar to what we have seen for $E_{6}$.

Proposition 13.30 (a) The number of positive roots of $E_{8}$ not in $E_{7}$ is 57.
(b) The subspace $V$ of $E_{8}$ spanned by vectors $e_{\alpha}$ for such roots is a 57-dimensional $E_{7}$-module. $V$ decomposes as the direct sum of a 56-dimensional fundamental module with a 1-dimensional module $L(0)$.

Proof. We see from Section 8.7 that the positive roots of $E_{8}$ not in $E_{7}$ are

$$
\begin{array}{ll}
\beta_{1}-\beta_{j} & j \in\{2,3,4,5,6,7\} \\
\beta_{1}+\beta_{j} & j \in\{2,3,4,5,6,7\} \\
\beta_{i}-\beta_{8} & i \in\{2,3,4,5,6,7\} \\
-\beta_{i}-\beta_{8} & i \in\{2,3,4,5,6,7\} \\
\beta_{1}-\beta_{8} & \\
\frac{1}{2} \sum_{i=1}^{8} \varepsilon_{i} \beta_{i} & \prod \varepsilon_{i}=1, \quad \varepsilon_{8}=-1, \quad \varepsilon_{1}=1
\end{array}
$$

The number of such roots is 57 .
Let $V$ be the subspace of $E_{8}$ spanned by the $e_{\alpha}$ for this set of roots. The argument of Proposition 13.29 shows that $V$ is an $E_{7}$-module. Now $\beta_{1}-\beta_{8}$ is orthogonal to all fundamental roots of $E_{7}$ and it follows that

$$
\left[e_{\alpha} e_{\beta_{1}-\beta_{8}}\right]=0 \quad \text { for all } \alpha \in \Phi\left(E_{7}\right)
$$

Hence $\mathbb{C} e_{\beta_{1}-\beta_{8}}$ is a 1-dimensional $E_{7}$-submodule of $V$. Let $V^{\prime}$ be the subspace spanned by the remaining $e_{\alpha}$. The fact that $\beta_{1}-\beta_{8}$ is orthogonal to all $\alpha \in \Phi\left(E_{7}\right)$ implies that $\beta_{1}-\beta_{8}$ cannot be expressed in the form $\alpha+\beta$ where $\alpha \in \Phi\left(E_{7}\right), \beta \in \Phi^{+}\left(E_{8}\right)$. This shows that $V^{\prime}$ is an $E_{7}$-submodule of $V$. Its highest weight is obtained by picking the weight with the highest value of $h$, and this is $\beta_{2}-\beta_{8}$. In fact the first few highest weights are

$$
\beta_{2}-\beta_{8}, \quad \beta_{3}-\beta_{8}, \quad \beta_{4}-\beta_{8}, \quad \beta_{5}-\beta_{8}, \quad \ldots
$$

By calculating scalar products $\{$,$\} with the fundamental roots of E_{7}$ we see that $\beta_{2}-\beta_{8}=\omega_{2}$. Thus $L\left(\omega_{2}\right)$ is an irreducible direct summand of $V^{\prime}$. Since

$$
\operatorname{dim} L\left(\omega_{2}\right)=56=\operatorname{dim} V^{\prime}
$$

we have $V^{\prime}=L\left(\omega_{2}\right)$. Thus

$$
V=L\left(\omega_{2}\right) \oplus L(0)
$$

We can obtain information about some of the other fundamental $E_{7}$-modules by considering exterior powers of $V^{\prime}$. The highest weight of $\Lambda^{2} V^{\prime}$ is

$$
\left(\beta_{2}-\beta_{8}\right)+\left(\beta_{3}-\beta_{8}\right)=\beta_{2}+\beta_{3}-2 \beta_{8}
$$

A calculation of scalar products $\{$,$\} shows that$

$$
\beta_{2}+\beta_{3}-2 \beta_{8}=\omega_{3}
$$

Thus $\Lambda^{2} V^{\prime}$ contains $L\left(\omega_{3}\right)$ as an irreducible direct summand. But we know that

$$
\operatorname{dim} L\left(\omega_{3}\right)=\binom{56}{2}-1
$$

Thus

$$
\Lambda^{2} V^{\prime}=L\left(\omega_{3}\right) \oplus L(0)
$$

The highest weight of $\Lambda^{3} V^{\prime}$ is

$$
\left(\beta_{2}-\beta_{8}\right)+\left(\beta_{3}-\beta_{8}\right)+\left(\beta_{4}-\beta_{8}\right)=\beta_{2}+\beta_{3}+\beta_{4}-3 \beta_{8}
$$

We have

$$
\beta_{2}+\beta_{3}+\beta_{4}-3 \beta_{8}=\omega_{4} .
$$

Thus $L\left(\omega_{4}\right)$ is an irreducible direct summand of $\Lambda^{3} V^{\prime}$. We know that

$$
\operatorname{dim} L\left(\omega_{4}\right)=\binom{56}{3}-56
$$

In fact we have

$$
\Lambda^{3} V^{\prime}=L\left(\omega_{4}\right) \oplus L\left(\omega_{2}\right)
$$

The highest weight of $\Lambda^{4} V^{\prime}$ is

$$
\left(\beta_{2}-\beta_{8}\right)+\left(\beta_{3}-\beta_{8}\right)+\left(\beta_{4}-\beta_{8}\right)+\left(\beta_{5}-\beta_{8}\right)=\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}-4 \beta_{8}
$$

We have

$$
\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}-4 \beta_{8}=\omega_{5} .
$$

We know that

$$
\operatorname{dim} L\left(\omega_{5}\right)=\binom{56}{4}-\binom{56}{2}
$$

In fact it turns out that

$$
\Lambda^{4} V^{\prime}=L\left(\omega_{5}\right) \oplus L\left(\omega_{3}\right) \oplus L(0)
$$

Some of the remaining fundamental $E_{7}$-modules may be identified by means of the adjoint module. The highest root of $E_{7}$ is $-\beta_{1}-\beta_{8}$ and we have $-\beta_{1}-\beta_{8}=\omega_{8}$. Thus we see that $L\left(\omega_{8}\right)$ is the adjoint $E_{7}$-module, since

$$
\operatorname{dim} L\left(\omega_{8}\right)=133=\operatorname{dim} L
$$

The second highest root of $E_{7}$ is $\frac{1}{2}\left(-\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}+\beta_{7}-\beta_{8}\right)$. Thus the highest weight of $\Lambda^{2} L$ is

$$
\begin{gathered}
\left(-\beta_{1}-\beta_{8}\right)+\frac{1}{2}\left(-\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}+\beta_{7}-\beta_{8}\right) \\
\quad=\frac{1}{2}\left(-3 \beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}+\beta_{7}-3 \beta_{8}\right) .
\end{gathered}
$$

We have

$$
\frac{1}{2}\left(-3 \beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}+\beta_{7}-3 \beta_{8}\right)=\omega_{7}
$$

Thus $L\left(\omega_{7}\right)$ is an irreducible direct summand of $\Lambda^{2} L$. Since

$$
\operatorname{dim} L\left(\omega_{7}\right)=\binom{133}{2}-1
$$

we have

$$
\Lambda^{2} L=L\left(\omega_{7}\right) \oplus L(0)
$$

We next consider the simple Lie algebra $E_{8}$. The smallest dimension of a fundamental module for $E_{8}$ is

$$
\operatorname{dim} L\left(\omega_{1}\right)=248
$$

The highest root of $E_{8}$ is $\beta_{1}-\beta_{8}$, and we have $\beta_{1}-\beta_{8}=\omega_{1}$. Since $\operatorname{dim} L=248$ we deduce that $L\left(\omega_{1}\right)=L$. Thus the fundamental module $L\left(\omega_{1}\right)$ is the adjoint module.

The description of the remaining fundamental modules of $E_{8}$ is considerably more complicated than in the other simple Lie algebras. We shall not discuss the details.

We now turn to the simple Lie algebra $F_{4}$ and show how to obtain the 26-dimensional fundamental module. This will be done by identifying $F_{4}$ with a subalgebra of $E_{6}$. We shall retain our previous numbering of the fundamental roots of $E_{6}$ given by


Let $\sigma$ be the permutation of the vertices given by

$$
\sigma=(3 \quad 8)(4 \quad 7)(5)(6)
$$

Then $\sigma$ gives a symmetry of the Dynkin diagram of $E_{6}$ with $\sigma^{2}=1$. We have

$$
A_{\sigma(i) \sigma(j)}=A_{i j} \quad \text { for all } i, j
$$

Thus by Theorem 7.5 there is an automorphism of $E_{6}$, which we shall also call $\sigma$, satisfying

$$
\begin{aligned}
& \sigma\left(e_{i}\right)=e_{\sigma(i)} \\
& \sigma\left(f_{i}\right)=f_{\sigma(i)} \\
& \sigma\left(h_{i}\right)=h_{\sigma(i)}
\end{aligned}
$$

Since $e_{i}, f_{i}, h_{i}$ generate the Lie algebra, $\sigma$ is determined by these conditions, and we have $\sigma^{2}=1$.

We may define a linear map on the real vector space spanned by the simple roots, also denoted by $\sigma$, to satisfy

$$
\sigma\left(\alpha_{i}\right)=\alpha_{\sigma(i)}
$$

Then we have $\sigma(\Phi)=\Phi$. All the $\sigma$-orbits on $\Phi$ have size 1 or 2 . Examination of the root system of $E_{6}$ shows there are 24 orbits of size 1 and 24 of size 2 .

Proposition 13.31 Let $L$ be the simple Lie algebra $E_{6}$ and $\sigma: L \rightarrow L$ be the automorphism of order 2 given above. Then the subalgebra $L^{\sigma}$ of $\sigma$-stable elements of $L$ is isomorphic to $F_{4}$. The elements

$$
\begin{array}{llll}
E_{1}=e_{6} & E_{2}=e_{5} & E_{3}=e_{4}+e_{7} & E_{4}=e_{3}+e_{8} \\
F_{1}=f_{6} & F_{2}=f_{5} & F_{3}=f_{4}+f_{7} & F_{4}=f_{3}+f_{8} \\
H_{1}=h_{6} & H_{2}=h_{5} & H_{3}=h_{4}+h_{7} & H_{4}=h_{3}+h_{8}
\end{array}
$$

are standard generators of $F_{4}$.

Proof. Let $\left(A_{i j}\right)$ be the Cartan matrix of $F_{4}$ given by

$$
A=\left(\begin{array}{cccc}
2 & -1 & 0 & 0 \\
-1 & 2 & -1 & 0 \\
0 & -2 & 2 & -1 \\
0 & 0 & -1 & 2
\end{array}\right)
$$

It is straightforward to check that the elements $E_{i}, F_{i}, H_{i}$ satisfy the relations

$$
\begin{aligned}
& {\left[H_{i} H_{j}\right]=0} \\
& {\left[H_{i} E_{j}\right]=A_{i j} E_{j}} \\
& {\left[H_{i} F_{j}\right]=-A_{i j} F_{j}}
\end{aligned}
$$

$$
\begin{aligned}
& {\left[E_{i} F_{i}\right]=H_{i}} \\
& {\left[E_{i} F_{j}\right]=0 \quad \text { if } i \neq j} \\
& {\left[E_{i}, \ldots\left[E_{i} E_{j}\right]\right]=0 \quad \text { if } i \neq j} \\
& {\left[F_{i}, \ldots\left[F_{i} F_{j}\right]\right]=0 \quad \text { if } i \neq j}
\end{aligned}
$$

where the last two relations have $1-A_{i j}$ factors $E_{i}, F_{i}$ respectively. By Proposition 7.35 there is a homomorphism

$$
\theta: F_{4} \rightarrow L
$$

whose image is the subalgebra generated by the elements $E_{i}, F_{i}, H_{i}$. Since $\theta \neq 0$ and $F_{4}$ is simple the image of $\theta$ is isomorphic to $F_{4}$.

We shall also show that $\operatorname{im} \theta=L^{\sigma}$. Since each $E_{i}, F_{i}, H_{i}$ lies in $L^{\sigma}$ we have $\operatorname{im} \theta \subset L^{\sigma}$. On the other hand consider the decomposition

$$
L=H \oplus \sum_{\substack{\alpha \in \Phi \\ \sigma(\alpha)=\alpha}} \mathbb{C} e_{\alpha} \oplus \sum_{\substack{\alpha \in \Phi \\ \sigma(\alpha) \neq \alpha}}\left(\mathbb{C} e_{\alpha}+\mathbb{C} e_{\sigma(\alpha)}\right)
$$

Each direct summand is $\sigma$-stable, thus $L^{\sigma}$ is the direct sum of the $\sigma$-stable subspaces of the components. We have

$$
\begin{aligned}
& \operatorname{dim} H^{\sigma}=4 \\
& \operatorname{dim}\left(\mathbb{C} e_{\alpha}\right)^{\sigma} \leq 1 \quad \text { if } \sigma(\alpha)=\alpha \\
& \operatorname{dim}\left(\mathbb{C} e_{\alpha}+\mathbb{C} e_{\sigma(\alpha)}\right)^{\sigma} \leq 1 \quad \text { if } \sigma(\alpha) \neq \alpha .
\end{aligned}
$$

Thus $\operatorname{dim} L^{\sigma} \leq 4+24+24=52$. But $\operatorname{dim}(\operatorname{im} \theta)=52$, thus $\operatorname{im} \theta=L^{\sigma}$. Hence $L^{\sigma}$ is isomorphic to $F_{4}$.

Now let $V$ be the 27-dimensional fundamental module $L\left(\omega_{8}\right)$ for $E_{6}$ constructed in Proposition 13.29. Then $V$ may be regarded as an $F_{4}$-module using our embedding of $F_{4}$ in $E_{6}$. We label the fundamental roots of $F_{4}$ by the diagram


Proposition 13.32 The $F_{4}$-module $V$ decomposes as

$$
V=L\left(\omega_{4}\right) \oplus L(0)
$$

where $L\left(\omega_{4}\right)$ is the 26-dimensional fundamental module.

Proof. We determine the weights of the $F_{4}$-module $V$. We recall that the weights of $V$ have form

$$
\begin{array}{ll}
\beta_{2}-\beta_{j} & 3 \leq j \leq 7 \\
\beta_{2}+\beta_{j} & 3 \leq j \leq 7 \\
-\beta_{1}-\beta_{8} & \\
\frac{1}{2} \sum_{i=1}^{8} \varepsilon_{i} \beta_{i} & \prod \varepsilon_{i}=1, \quad \varepsilon_{1}=-1, \quad \varepsilon_{2}=1, \quad \varepsilon_{8}=-1
\end{array}
$$

Now the fundamental roots of $E_{8}$ are

$$
\begin{aligned}
& \alpha_{i}=\beta_{i}-\beta_{i+1} \quad i=1, \ldots, 6 \\
& \alpha_{7}=\beta_{6}+\beta_{7} \\
& \alpha_{8}=-\frac{1}{2} \sum_{i=1}^{8} \beta_{i} .
\end{aligned}
$$

Also $\alpha_{j}\left(h_{i}\right)=A_{i j} \quad i, j \in\{1, \ldots, 8\}$ where $\left(A_{i j}\right)$ is the Cartan matrix of $E_{8}$. It follows that the numbers $\beta_{j}\left(h_{i}\right) \quad i, j \in\{1, \ldots, 8\}$ are given by

$$
\begin{array}{lc}
\beta_{i}\left(h_{i}\right)=1 & i=1, \ldots, 7 \\
\beta_{i+1}\left(h_{i}\right)=-1 & i=1, \ldots, 6 \\
\beta_{i}\left(h_{8}\right)=-\frac{1}{2} & i=1, \ldots, 8 \\
\beta_{i}\left(h_{j}\right)=0 & \text { otherwise } .
\end{array}
$$

Let $H_{1}, H_{2}, H_{3}, H_{4}$ be the fundamental coroots of $F_{4}$ defined above and $\omega_{1}, \omega_{2}, \omega_{3}, \omega_{4}$ the corresponding fundamental weights of $F_{4}$. Then $\omega_{i}\left(H_{j}\right)=\delta_{i j}$. By calculating the values $\beta_{i}\left(H_{j}\right)$ we deduce

$$
\begin{array}{ll}
\beta_{1}=\beta_{2}=\beta_{8} & =-\frac{1}{2} \omega_{4} \\
\beta_{3} & =\frac{1}{2} \omega_{4} \\
\beta_{4} & =\omega_{3}-\frac{3}{2} \omega_{4} \\
\beta_{5} & =\omega_{2}-\omega_{3}-\frac{1}{2} \omega_{4} \\
\beta_{6} & =\omega_{1}-\omega_{2}+\omega_{3}-\frac{1}{2} \omega_{4} \\
\beta_{7} & =-\omega_{1}+\omega_{3}-\frac{1}{2} \omega_{4}
\end{array}
$$

when the $\beta_{i}$ are regarded as weights for $F_{4}$. Hence the 27 weights of the $F_{4}$-module $V$ are

$$
\begin{aligned}
& \pm\left\{\omega_{4}, \omega_{1}-\omega_{3}, \omega_{1}-\omega_{4}, \omega_{2}-\omega_{3}, \omega_{3}-\omega_{4}, \omega_{3}-2 \omega_{4}, \omega_{1}-\omega_{2}+\omega_{3}, \omega_{1}-\omega_{2}\right. \\
& \left.\quad+\omega_{4}, \omega_{1}-\omega_{3}+\omega_{4}, \omega_{2}-\omega_{3}-\omega_{4}, \omega_{2}-2 \omega_{3}+\omega_{4}, \omega_{1}-\omega_{2}+\omega_{3}-\omega_{4}\right\} \\
& \cup\{0,0,0\}
\end{aligned}
$$

The only dominant weight among these, excluding 0 , is $\omega_{4}$. Thus $V$ has highest weight $\omega_{4}$ and so $L\left(\omega_{4}\right)$ is an irreducible direct summand of $V$. Since

$$
\operatorname{dim} V=27, \quad \operatorname{dim} L\left(\omega_{4}\right)=26
$$

we have

$$
V=L\left(\omega_{4}\right) \oplus L(0)
$$

Using the relation $\alpha_{i}=\sum_{j} A_{j i} \omega_{j}$ in $F_{4}$ we see that

$$
\omega_{4}=\alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+2 \alpha_{4} .
$$

This is the highest short root of $F_{4}$. All short roots of $F_{4}$ are transforms of this one under elements of the Weyl group $W$. Thus all 24 short roots of $F_{4}$ are weights of $L\left(\omega_{4}\right)$. So the weights of $L\left(\omega_{4}\right)$ are the 24 short roots together with 0 with multiplicity 2 .

We now discuss the other fundamental modules for $F_{4}$. We first consider $L\left(\omega_{1}\right)$. The relations $\alpha_{i}=\sum_{j} A_{j i} \omega_{j}$ for $F_{4}$ show that

$$
\omega_{1}=2 \alpha_{1}+3 \alpha_{2}+4 \alpha_{3}+2 \alpha_{4}
$$

We recall from Section 8.6 that

$$
\alpha_{1}=\beta_{1}-\beta_{2} \quad \alpha_{2}=\beta_{2}-\beta_{3} \quad \alpha_{3}=\beta_{3} \quad \alpha_{4}=\frac{1}{2}\left(-\beta_{1}-\beta_{2}-\beta_{3}+\beta_{4}\right)
$$

and so

$$
\omega_{1}=2 \alpha_{1}+3 \alpha_{2}+4 \alpha_{3}+2 \alpha_{4}=\beta_{1}+\beta_{4}
$$

The long roots of $F_{4}$ have form $\pm \beta_{i} \pm \beta_{j}$ and, since $\beta_{4} \succ \beta_{1} \succ \beta_{2} \succ \beta_{3}, \beta_{1}+\beta_{4}$ is the highest root. Thus $\omega_{1}$ is the highest root of $F_{4}$ and $L\left(\omega_{1}\right)$ is therefore the adjoint $F_{4}$-module.

The remaining fundamental modules $L\left(\omega_{2}\right), L\left(\omega_{3}\right)$ for $F_{4}$ satisfy

$$
\begin{aligned}
& \operatorname{dim} L\left(\omega_{2}\right)=\binom{52}{2}-52 \\
& \operatorname{dim} L\left(\omega_{3}\right)=\binom{26}{2}-52
\end{aligned}
$$

It can be shown that $L\left(\omega_{2}\right), L\left(\omega_{3}\right)$ appear as irreducible direct summands of $\Lambda^{2} L\left(\omega_{1}\right), \Lambda^{2} L\left(\omega_{4}\right)$ respectively, and that

$$
\begin{aligned}
& \Lambda^{2} L\left(\omega_{1}\right)=L\left(\omega_{1}\right) \oplus L\left(\omega_{2}\right) \\
& \Lambda^{2} L\left(\omega_{4}\right)=L\left(\omega_{1}\right) \oplus L\left(\omega_{3}\right) .
\end{aligned}
$$

Finally we consider the simple Lie algebra $G_{2}$ and show how to obtain the 7-dimensional fundamental module. We do this by identifying $G_{2}$ with a subalgebra of $D_{4}$. The fundamental roots of $D_{4}$ will be numbered as in the diagram


Let $\sigma$ be the permutation of the vertices given by

$$
\sigma=\left(\begin{array}{lll}
1 & 3 & 4
\end{array}\right)(2) .
$$

$\sigma$ gives a symmetry of the Dynkin diagram with $\sigma^{3}=1$. Since

$$
A_{\sigma(i) \sigma(j)}=A_{i j} \quad \text { for all } i, j
$$

there exists by Theorem 7.5 an automorphism $\sigma$ of $D_{4}$ satisfying

$$
\begin{aligned}
& \sigma\left(e_{i}\right)=e_{\sigma(i)} \\
& \sigma\left(f_{i}\right)=f_{\sigma(i)} \\
& \sigma\left(h_{i}\right)=h_{\sigma(i)}
\end{aligned}
$$

We may also define a linear map $\sigma$ on the vector space spanned by the simple roots, satisfying $\sigma\left(\alpha_{i}\right)=\alpha_{\sigma(i)}$. We have $\sigma(\Phi)=\Phi$. These are $6 \sigma$-orbits of size 1 on $\Phi$ and 6 orbits of size 3 .

Proposition 13.33 Let L be the simple Lie algebra $D_{4}$ and $\sigma: L \rightarrow L$ be the automorphism of order 3 given above. Then the subalgebra $L^{\sigma}$ of $\sigma$-stable elements of $L$ is isomorphic to $G_{2}$.

The elements

$$
\begin{array}{ll}
E_{1}=e_{2} & E_{2}=e_{1}+e_{3}+e_{4} \\
F_{1}=f_{2} & F_{2}=f_{1}+f_{3}+f_{4} \\
H_{1}=h_{2} & H_{2}=h_{1}+h_{3}+h_{4}
\end{array}
$$

are standard generators of $G_{2}$.

Proof. The idea is the same as that for $F_{4}$ in $E_{6}$. The Cartan matrix of $G_{2}$ is

$$
A=\left(\begin{array}{cc}
2 & -1 \\
-3 & 2
\end{array}\right)
$$

It is again straightforward to check that the elements $E_{1}, E_{2}, F_{1}, F_{2}, H_{1}, H_{2}$ satisfy the defining relations

$$
\begin{aligned}
& {\left[H_{i} H_{j}\right]=0} \\
& {\left[H_{i} E_{j}\right]=A_{i j} E_{j}} \\
& {\left[H_{i} F_{j}\right]=-A_{i j} F_{j}} \\
& {\left[E_{i} F_{i}\right]=H_{i}} \\
& {\left[E_{i} F_{j}\right]=0 \quad \text { if } i \neq j} \\
& {\left[E_{i}, \ldots\left[E_{i} E_{j}\right]\right]=0 \quad \text { if } i \neq j} \\
& {\left[F_{i}, \ldots\left[F_{i} F_{j}\right]\right]=0 \quad \text { if } i \neq j}
\end{aligned}
$$

where the last two relations have $1-A_{i j}$ factors $E_{i}, F_{i}$ respectively.
Thus by Proposition 7.35 there is a homomorphism $\theta: G_{2} \rightarrow L$. The image $\operatorname{im} \theta$ is isomorphic to $G_{2}$. We show $\operatorname{im} \theta=L^{\sigma}$. Since $E_{i}, F_{i}, H_{i}$ lie in $L^{\sigma}$ we have $\operatorname{im} \theta \subset L^{\sigma}$. Now consider the decomposition

$$
L=H \oplus \sum_{\substack{\alpha \in \Phi \\ \sigma(\alpha)=\alpha}} \mathbb{C} e_{\alpha} \oplus \sum_{\substack{\alpha \in \Phi \\ \sigma(\alpha) \neq \alpha}}\left(\mathbb{C} e_{\alpha}+\mathbb{C} e_{\sigma(\alpha)}+\mathbb{C} e_{\sigma^{2}(\alpha)}\right)
$$

Each direct summand is $\sigma$-stable, thus $L^{\sigma}$ is the direct sum of the $\sigma$-stable subspaces of the components. We have

$$
\begin{aligned}
& \operatorname{dim} H^{\sigma}=2 \\
& \operatorname{dim}\left(\mathbb{C} e_{\alpha}\right)^{\sigma} \leq 1 \quad \text { if } \sigma(\alpha)=\alpha \\
& \operatorname{dim}\left(\mathbb{C} e_{\alpha}+\mathbb{C} e_{\sigma(\alpha)}+\mathbb{C} e_{\sigma^{2}(\alpha)}\right)^{\sigma} \leq 1 \quad \text { if } \sigma(\alpha) \neq \alpha
\end{aligned}
$$

Thus $\operatorname{dim} L^{\sigma} \leq 2+6+6=14$. But $\operatorname{dim}(\operatorname{im} \theta)=14$, thus $\operatorname{im} \theta=L^{\sigma}$. Hence $L^{\sigma}$ is isomorphic to $G_{2}$.

Proposition 13.34 Let $V$ be the 8 -dimensional natural $D_{4}$-module. Regard $V$ as a $G_{2}$-module using the above embedding of $G_{2}$ in $D_{4}$. Then

$$
V=L\left(\omega_{2}\right) \oplus L(0)
$$

where $L\left(\omega_{2}\right)$ is the 7-dimensional fundamental $G_{2}$-module.

Proof. We recall from Section 8.2 that in this 8-dimensional representation we have
$e_{1}=E_{12}-E_{-2-1}, \quad e_{2}=E_{23}-E_{-3-2}, \quad e_{3}=E_{34}-E_{-4-3}, \quad e_{4}=E_{3-4}-E_{4-3}$
$f_{1}=-E_{-1-2}+E_{21}, \quad f_{2}=-E_{-2-3}+E_{32}, \quad f_{3}=-E_{-3-4}+E_{43}$,
$f_{4}=-E_{-34}+E_{-43}$.
Hence

$$
\begin{aligned}
& h_{1}=E_{11}-E_{22}-E_{-1-1}+E_{-2-2} \\
& h_{2}=E_{22}-E_{33}-E_{-2-2}+E_{-3-3} \\
& h_{3}=E_{33}-E_{44}-E_{-3-3}+E_{-4-4} \\
& h_{4}=E_{33}-E_{-4-4}-E_{-3-3}+E_{44}
\end{aligned}
$$

and so

$$
\begin{aligned}
& H_{1}=E_{22}-E_{33}-E_{-2-2}+E_{-3-3} \\
& H_{2}=E_{11}-E_{22}+2 E_{33}-E_{-1-1}+E_{-2-2}-2 E_{-3-3} .
\end{aligned}
$$

Let $v_{1}, v_{2}, v_{3}, v_{4}, v_{-1}, v_{-2}, v_{-3}, v_{-4}$ be the natural basis of $V$. Let $\omega_{1}, \omega_{2}$ be the fundamental weights for $G_{2}$. Since $\omega_{i}\left(H_{j}\right)=\delta_{i j}$ these basis vectors span weight spaces with weights

$$
\omega_{2}, \quad \omega_{1}-\omega_{2}, \quad-\omega_{1}+2 \omega_{2}, \quad 0, \quad-\omega_{2}, \quad-\omega_{1}+\omega_{2}, \quad \omega_{1}-2 \omega_{2}, \quad 0
$$

respectively. The highest weight is $\omega_{2}$, thus $L\left(\omega_{2}\right)$ is an irreducible direct summand of $V$. We have

$$
\operatorname{dim} V=8, \quad \operatorname{dim} L\left(\omega_{2}\right)=7
$$

and so

$$
V=L\left(\omega_{2}\right) \oplus L(0)
$$

We note that $\omega_{2}=\alpha_{1}+2 \alpha_{2}$ is the highest short root of $G_{2}$. All short roots are transforms of this root by elements of the Weyl group, thus all six short roots are weights of $L\left(\omega_{2}\right)$. Thus the weights of $L\left(\omega_{2}\right)$ are the short roots together with 0 .

Now we have

$$
\begin{aligned}
& E_{1}=E_{23}-E_{-3-2}, \quad E_{2}=E_{12}+E_{34}+E_{3-4}-E_{-2-1}-E_{-4-3}-E_{4-3} \\
& F_{1}=-E_{-2-3}+E_{32}, \quad F_{2}=-E_{-1-2}-E_{-3-4}-E_{-34}+E_{21}+E_{43}+E_{-43} .
\end{aligned}
$$

It may be checked that the vector $v_{4}-v_{-4}$ is annihilated by $E_{1}, E_{2}, F_{1}, F_{2}$ and so spans the 1 -dimensional submodule $L(0)$.

Finally we consider the other fundamental $G_{2}$-module $L\left(\omega_{1}\right)$. The relations $\alpha_{i}=\sum_{j} A_{j i} \omega_{j}$ show that

$$
\omega_{1}=2 \alpha_{1}+3 \alpha_{2}
$$

This is the highest root of $G_{2}$. Therefore the fundamental module $L\left(\omega_{1}\right)$ is the 14 -dimensional adjoint module.

## 14

## Generalised Cartan matrices and Kac-Moody algebras

In 1967 V. G. Kac and R. V. Moody independently initiated the study of certain Lie algebras $L(A)$ associated with a generalised Cartan matrix $A$. An $n \times n$ matrix $A=\left(A_{i j}\right)$ is called a generalised Cartan matrix if it satisfies the conditions

$$
\begin{aligned}
& A_{i i}=2 \quad \text { for } i=1, \ldots, n \\
& A_{i j} \in \mathbb{Z} \quad \text { and } \quad A_{i j} \leq 0 \quad \text { if } i \neq j \\
& A_{i j}=0 \quad \text { implies } \quad A_{j i}=0 .
\end{aligned}
$$

The Cartan matrix of any finite dimensional simple Lie algebra is a generalised Cartan matrix, as shown in Section 6.4. We shall see that, in the special case when $A$ is a Cartan matrix, the Lie algebra $L(A)$ constructed by Kac and Moody coincides with the finite dimensional simple Lie algebra with Cartan matrix $A$. However, the Lie algebra $L(A)$ can in general be infinite dimensional.

The term 'generalised Cartan matrix' will be abbreviated to GCM. The Lie algebra $L(A)$ associated to a GCM $A$ will be called the Kac-Moody algebra associated to $A$. We shall explain the definition and some of the basic properties of $L(A)$ in the present chapter. In fact the introductory ideas do not use the fact that $A$ is a GCM - we shall assume initially that $A$ is any $n \times n$ matrix over $\mathbb{C}$.

### 14.1 Realisations of a square matrix

Let $A$ be an $n \times n$ matrix over $\mathbb{C}$. A realisation of $A$ is a triple $\left(H, \Pi, \Pi^{\mathrm{v}}\right)$ where:
$H$ is a finite dimensional vector space over $\mathbb{C}$
$\Pi^{\mathrm{v}}=\left\{h_{1}, \ldots, h_{n}\right\}$ is a linearly independent subset of $H$
$\Pi=\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ is a linearly independent subset of $H^{*}$ $\alpha_{j}\left(h_{i}\right)=A_{i j}$ for all $i, j$.

Proposition 14.1 If $\left(H, \Pi, \Pi^{\mathrm{v}}\right)$ is a realisation of $A$ then $\operatorname{dim} H \geq 2 n-$ rank $A$.

Proof. Let rank $A=l$ and $\operatorname{dim} H=m$. We extend the set $\Pi^{\mathrm{v}}$ to give a basis $h_{1}, \ldots, h_{m}$ of $H$ and extend $\Pi$ to give a basis $\alpha_{1}, \ldots, \alpha_{m}$ of $H^{*}$. Consider the $m \times m$ matrix $\left(\alpha_{j}\left(h_{i}\right)\right)$. This is non-singular so its rows are linearly independent. Thus the $n \times m$ matrix given by the first $n$ rows has rank $n$. This matrix therefore has $n$ linearly independent columns. Now the leading $n \times n$ submatrix is $A$, so has rank $l$. Thus the remaining $n \times(m-n)$ matrix has rank at least $n-l$. It follows that $m-n \geq n-l$, that is $m \geq 2 n-l$.

Definition A minimal realisation of $A$ is a realisation in which

$$
\operatorname{dim} H=2 n-\operatorname{rank} A
$$

Proposition 14.2 Any $n \times n$ matrix over $\mathbb{C}$ has a minimal realisation.

Proof. Since rank $A=l, A$ has a non-singular $l \times l$ submatrix. By reordering the rows and columns we obtain a matrix

$$
\begin{array}{cc}
l \\
n-l
\end{array}\left(\begin{array}{cc}
A_{11} & A_{12} \\
A_{21} & A_{22}
\end{array}\right)
$$

in which $A_{11}$ is non-singular. Let

$$
\begin{gathered}
C=\left(\begin{array}{ccc}
A_{11} & A_{12} & O \\
A_{21} & A_{22} & I_{n-l} \\
O & I_{n-l} & O
\end{array}\right) \begin{array}{c}
l \\
n-l \\
n-l
\end{array} \\
l \begin{array}{lll}
n-l & n-l
\end{array}
\end{gathered}
$$

Since $\operatorname{det} C= \pm \operatorname{det} A_{11} \neq 0$ we see that $C$ is a non-singular $(2 n-l) \times(2 n-l)$ matrix. Let $H$ be the vector space of all $(2 n-l)$-tuples over $\mathbb{C}$. Define $\alpha_{1}, \ldots, \alpha_{n} \in H^{*}$ to be the first $n$ coordinate functions

$$
\left(\lambda_{1}, \ldots, \lambda_{2 n-l}\right) \rightarrow \lambda_{i} \quad i=1, \ldots, n
$$

Define $h_{1}, \ldots, h_{n} \in H$ to be the first $n$ row vectors of $C$. Then $\alpha_{1}, \ldots, \alpha_{n}$ and $h_{1}, \ldots, h_{n}$ are linearly independent and we obtain a realisation of

$$
\left(\begin{array}{ll}
A_{11} & A_{12} \\
A_{21} & A_{22}
\end{array}\right)
$$

with $\operatorname{dim} H=2 n-l$. By reordering $\alpha_{1}, \ldots, \alpha_{n}$ and $h_{1}, \ldots, h_{n}$ appropriately we obtain a minimal realisation of $A$.

Now let $\left(H, \Pi, \Pi^{\mathrm{v}}\right)$ and $\left(H^{\prime}, \Pi^{\prime},\left(\Pi^{\prime}\right)^{\mathrm{v}}\right)$ be two realisations of $A$. We say the realisations are isomorphic if there is an isomorphism of vector spaces

$$
\phi: H \rightarrow H^{\prime}
$$

such that $\phi\left(h_{i}\right)=h_{i}^{\prime}$ and $\phi^{*}\left(\alpha_{i}^{\prime}\right)=\alpha_{i}$ where

$$
\phi^{*}:\left(H^{\prime}\right)^{*} \rightarrow H^{*}
$$

is the isomorphism induced by $\phi$.

Proposition 14.3 Any two minimal realisations of an $n \times n$ matrix $A$ over $\mathbb{C}$ are isomorphic.

Proof. Let $\left(H, \Pi, \Pi^{\mathrm{v}}\right)$ be the minimal realisation of $A$ constructed in Proposition 14.2 and $\left(H^{\prime}, \Pi^{\prime},\left(\Pi^{\prime}\right)^{\mathrm{v}}\right)$ be another minimal realisation. We reorder the rows and columns of $A$ as before to obtain

$$
\left(\begin{array}{ll}
A_{11} & A_{12} \\
A_{21} & A_{22}
\end{array}\right)
$$

where $A_{11}$ is non-singular.
We complete $h_{1}^{\prime}, \ldots, h_{n}^{\prime}$ to a basis $h_{1}^{\prime}, \ldots, h_{2 n-l}^{\prime}$ of $H^{\prime}$. Then the matrix $\left(\alpha_{j}^{\prime}\left(h_{i}^{\prime}\right)\right)$ for $i=1, \ldots, 2 n-l ; j=1, \ldots, n$ has form

$$
\left(\begin{array}{cc}
A_{11} & A_{12} \\
A_{21} & A_{22} \\
B_{1} & B_{2}
\end{array}\right)
$$

Since $\alpha_{1}^{\prime}, \ldots, \alpha_{n}^{\prime}$ are linearly independent this matrix has rank $n$. Thus it has $n$ linearly independent rows. Since rows $l+1, \ldots, n$ are linear combinations of rows $1, \ldots, l$ the matrix

$$
\begin{gathered}
\left(\begin{array}{cc}
A_{11} & A_{12} \\
B_{1} & B_{2}
\end{array}\right)_{n-l}^{l} \\
l \\
l
\end{gathered}
$$

must have linearly independent rows, so is non-singular.

We now extend $\alpha_{1}^{\prime}, \ldots, \alpha_{n}^{\prime}$ to $\alpha_{1}^{\prime}, \ldots, \alpha_{2 n-l}^{\prime}$ so that the $(2 n-l) \times(2 n-l)$ matrix $\left(\alpha_{j}^{\prime}\left(h_{i}^{\prime}\right)\right)$ is

$$
\begin{array}{ccc}
\left(\begin{array}{ccc}
A_{11} & A_{12} & O \\
A_{21} & A_{22} & I_{n-l} \\
B_{1} & B_{2} & O
\end{array}\right) \begin{array}{c}
l \\
n-l \\
l
\end{array} \quad n-l & n-l
\end{array}
$$

This matrix is non-singular, thus $\alpha_{1}^{\prime}, \ldots, \alpha_{2 n-l}^{\prime}$ are a basis for $\left(H^{\prime}\right)^{*}$.
Since $A_{11}$ is non-singular, by adding suitable linear combinations of the first $l$ rows to the last $n-l$ rows we may achieve $B_{1}=O$. Thus it is possible to choose $h_{n+1}^{\prime}, \ldots, h_{2 n-l}^{\prime}$ so that $h_{1}^{\prime}, \ldots, h_{2 n-l}^{\prime}$ are a basis of $H^{\prime}$ and

$$
\left(\alpha_{j}^{\prime}\left(h_{i}^{\prime}\right)\right)=\left(\begin{array}{ccc}
A_{11} & A_{12} & O \\
A_{21} & A_{22} & I_{n-l} \\
O & B_{2}^{\prime} & O
\end{array}\right)
$$

The matrix $B_{2}^{\prime}$ must be non-singular since the whole matrix is non-singular.
We now make a further change to $h_{n+1}^{\prime}, \ldots, h_{2 n-l}^{\prime}$ equivalent to left multiplying the above matrix by

$$
\left(\begin{array}{ccc}
I_{l} & O & O \\
O & I_{n-l} & O \\
O & O & \left(B_{2}^{\prime}\right)^{-1}
\end{array}\right)
$$

Then we obtain

$$
\left(\alpha_{j}^{\prime}\left(h_{i}^{\prime}\right)\right)=\left(\begin{array}{ccc}
A_{11} & A_{12} & O \\
A_{21} & A_{22} & I_{n-l} \\
O & I_{n-l} & O
\end{array}\right)
$$

This is equal to the matrix $C$ above. Thus the map $h_{i} \rightarrow h_{i}^{\prime}$ gives an isomorphism $H \rightarrow H^{\prime}$ which induces the isomorphism $\left(H^{\prime}\right)^{*} \rightarrow H^{*}$ given by $\alpha_{j}^{\prime} \rightarrow \alpha_{j}$. This shows that the realisations $\left(H, \Pi, \Pi^{v}\right)$ and $\left(H^{\prime}, \Pi^{\prime},\left(\Pi^{\prime}\right)^{\mathrm{v}}\right)$ are isomorphic.

### 14.2 The Lie algebra $\tilde{L}(A)$ associated with a complex matrix

Let $A$ be an $n \times n$ matrix over $\mathbb{C}$ with rank $l$. Let $\left(H, \Pi, \Pi^{v}\right)$ be a minimal realisation of $A$. Then we have

$$
\begin{aligned}
& \operatorname{dim} H=2 n-l \\
& \Pi^{\mathrm{v}}=\left\{h_{1}, \ldots, h_{n}\right\} \subset H, \quad \Pi=\left\{\alpha_{1}, \ldots, \alpha_{n}\right\} \subset H^{*} \\
& \alpha_{j}\left(h_{i}\right)=A_{i j}
\end{aligned}
$$

We define a Lie algebra $\tilde{L}(A)$ by generators and relations.
Let $X=\left\{e_{1}, \ldots, e_{n}, f_{1}, \ldots, f_{n}, \tilde{x}\right.$ for all $\left.x \in H\right\}$ and let $R$ be the following set of Lie words in $X$ :

$$
\begin{aligned}
& \tilde{x}-\lambda \tilde{y}-\mu \tilde{z} \quad \text { for all } x, y, z \in H, \quad \lambda, \mu \in \mathbb{C} \text { with } x=\lambda y+\mu z \\
& {[\tilde{x} \tilde{y}] \quad \text { for all } x, y \in H} \\
& {\left[e_{i} f_{i}\right]-\tilde{h}_{i} \quad \text { for } i=1, \ldots, n} \\
& {\left[e_{i} f_{j}\right] \quad \text { for all } i \neq j} \\
& {\left[\tilde{x} e_{i}\right]-\alpha_{i}(x) e_{i} \quad \text { for all } x \in H \text { and } i=1, \ldots, n} \\
& {\left[\tilde{x} f_{i}\right]+\alpha_{i}(x) f_{i} \quad \text { for all } x \in H \text { and } i=1, \ldots, n .}
\end{aligned}
$$

We define $\tilde{L}(A)=L(X ; R)$ to be the Lie algebra generated by the elements $X$ subject to relations $R$.

Lemma 14.4 If a different minimal realisation of $A$ is chosen the Lie algebra $\tilde{L}(A)$ is the same up to isomorphism.

Proof. This follows from Proposition 14.3.
We note that if $A$ is a Cartan matrix then $\tilde{L}(A)$ is the Lie algebra investigated earlier in Section 7.4 and Example 9.13. For in this case $A$ is non-singular and $H$ is the vector space with basis $h_{i}=\left[e_{i} f_{i}\right]$.

Proposition 14.5 There is an automorphism $\tilde{\omega}$ of $\tilde{L}(A)$ uniquely determined by

$$
\tilde{\omega}\left(e_{i}\right)=-f_{i}, \quad \tilde{\omega}\left(f_{i}\right)=-e_{i}, \quad \tilde{\omega}(\tilde{x})=-\tilde{x}
$$

for all $x \in H$. Also $\tilde{\omega}^{2}=1$.

Proof. There is a map $\tilde{\omega}: X \rightarrow F L(X)$ given by the above formulae. By Proposition 9.9 there is a unique Lie algebra homomorphism $F L(X) \rightarrow F L(X)$ extending this map. We shall denote this map also by $\tilde{\omega}$. It satisfies $\tilde{\omega}^{2}=1$. Let $\langle R\rangle$ be the ideal of $F L(X)$ generated by the above set $R$ of Lie words. By applying $\tilde{\omega}$ to the elements of $R$ we see that $\tilde{\omega}(\langle R\rangle) \subset\langle R\rangle$. Thus we may define the induced map

$$
\tilde{\omega}: F L(X) /\langle R\rangle \rightarrow F L(X) /\langle R\rangle .
$$

Since $\tilde{\omega}^{2}=1, \tilde{\omega}$ is an automorphism of $\tilde{L}(A)$.

Let $\tilde{H}$ be the subalgebra of $\tilde{L}(A)$ generated by the elements $\tilde{x}$ for all $x \in H$. Let $\tilde{N}$ be the subalgebra generated by $e_{1}, \ldots, e_{n}$ and $\tilde{N}^{-}$the subalgebra generated by $f_{1}, \ldots, f_{n}$. Then we have

$$
\tilde{\omega}(\tilde{H})=\tilde{H}, \quad \tilde{\omega}(\tilde{N})=\tilde{N}^{-}, \quad \tilde{\omega}\left(\tilde{N}^{-}\right)=\tilde{N}
$$

Now let $V$ be an $n$-dimensional vector space over $\mathbb{C}$ with basis $v_{1}, \ldots, v_{n}$ and let

$$
T(V)=\bigoplus_{s \geq 0} T^{s}(V)
$$

be the tensor algebra of $V$. Thus $T^{s}(V)$ has basis

$$
v_{i_{1}} \otimes \cdots \otimes v_{i_{s}}=v_{i_{1}} \ldots v_{i_{s}}
$$

for all $i_{1}, \ldots, i_{s} \in\{1, \ldots, n\}$. For each linear map $\lambda \in H^{*}$ we define a map

$$
\theta_{\lambda}: X \rightarrow \text { End } T(V)
$$

It is sufficient to define the effect of these endomorphisms on the basis elements of $T(V) . T^{0}(V)$ has basis 1 . We define

$$
\begin{aligned}
& \theta_{\lambda}(\tilde{x}) \cdot 1=\lambda(x) 1 \\
& \theta_{\lambda}(\tilde{x}) \cdot\left(v_{i_{1}} \ldots v_{i_{s}}\right)=\left(\lambda-\alpha_{i_{1}}-\cdots-\alpha_{i_{s}}\right)(x) v_{i_{1}} \ldots v_{i_{s}}
\end{aligned}
$$

for $x \in H$.

$$
\begin{aligned}
& \theta_{\lambda}\left(f_{j}\right) \cdot 1=v_{j} \\
& \theta_{\lambda}\left(f_{j}\right) \cdot\left(v_{i_{1}} \ldots v_{i_{s}}\right)=v_{j} v_{i_{1}} \ldots v_{i_{s}} .
\end{aligned}
$$

We define $\theta_{\lambda}\left(e_{j}\right)$ by induction on $s$ as follows

$$
\begin{aligned}
& \theta_{\lambda}\left(e_{j}\right) \cdot 1=0 \\
& \theta_{\lambda}\left(e_{j}\right) \cdot v_{i}=\delta_{i j} \lambda\left(h_{j}\right) 1 \\
& \theta_{\lambda}\left(e_{j}\right) \cdot\left(v_{i_{1}} \ldots v_{i_{s}}\right)=v_{i_{1}}\left(\theta_{\lambda}\left(e_{j}\right)\left(v_{i_{2}} \ldots v_{i_{s}}\right)\right) \\
& +\delta_{i j}\left(\lambda-\alpha_{i_{2}}-\cdots-\alpha_{i_{s}}\right)\left(h_{j}\right) v_{i_{2}} \ldots v_{i_{s}} \quad s>1 .
\end{aligned}
$$

Proposition 14.6 The above map $\theta_{\lambda}: X \rightarrow$ End $T(V)$ can be extended to a Lie algebra homomorphism $\tilde{L}(A) \rightarrow[$ End $T(V)]$.

Proof. The idea of the proof is essentially the same as in Proposition 7.9. $\theta_{\lambda}$ can first be extended to a homomorphism

$$
\theta_{\lambda}: F L(X) \rightarrow[\text { End } T(V)]
$$

by Proposition 9.9. We have

$$
\tilde{L}(A) \cong F L(X) /\langle R\rangle
$$

and so in order to show that $\theta_{\lambda}$ induces a homomorphism $\tilde{L}(A) \rightarrow[$ End $T(V)]$ we must verify that $\theta_{\lambda}(r)=0$ for all $r \in R$.

The elements of $R$ have form

$$
\begin{aligned}
& \tilde{x}-\lambda \tilde{y}-\mu \tilde{z} \\
& {[\tilde{x} \tilde{y}]} \\
& {\left[e_{i} f_{i}\right]-\tilde{h}_{i}} \\
& {\left[e_{i} f_{j}\right] \quad i \neq j} \\
& {\left[\tilde{x} e_{i}\right]-\alpha_{i}(x) e_{i}} \\
& {\left[\tilde{x} f_{i}\right]+\alpha_{i}(x) f_{i} .}
\end{aligned}
$$

The relation $\theta_{\lambda}(r)=0$ may be checked for each such $r \in R$ in a straightforward manner, just as in the proof of Proposition 7.9

Corollary 14.7 The map $x \rightarrow \tilde{x}$ is an isomorphism of vector spaces $H \rightarrow \tilde{H}$.

Proof. $\tilde{H}$ is the subalgebra of $\tilde{L}(A)$ generated by $\tilde{x}$ for all $x \in H$. However, these elements form a Lie algebra since

$$
\begin{aligned}
& \tilde{x}_{1}+\tilde{x}_{2}=x_{1} \widetilde{+} x_{2} \\
& \lambda \tilde{x}=\widetilde{\lambda x} \\
& {\left[\tilde{x}_{1} \tilde{x}_{2}\right]=0 .}
\end{aligned}
$$

Thus $\tilde{H}=\{\tilde{x} ; \quad x \in H\}$.
Consider the map $H \rightarrow \tilde{H}$ given by $x \rightarrow \tilde{x}$. This is a homomorphism of Lie algebras. It is surjective. To show it is an isomorphism we must show it is also injective. Thus suppose $x \in H$ and $\tilde{x}=0$. Then $\theta_{\lambda}(\tilde{x})=0$. Thus $\lambda(x)=0$. Since this holds for all $\lambda \in H^{*}$ we may deduce that $x=0$.

We next consider the restriction of $\theta_{\lambda}$ to $\tilde{N}^{-}$. It is clear from the definition that this is independent of $\lambda$. We call it

$$
\theta: \tilde{N}^{-} \rightarrow[\text { End } T(V)]
$$

Now $\theta\left(f_{i}\right)$ is left multiplication by $v_{i}$. Thus, for any Lie word $w\left(f_{1}, \ldots, f_{n}\right)$ in $f_{1}, \ldots, f_{n}, \theta\left(w\left(f_{1}, \ldots, f_{n}\right)\right)$ is left multiplication by $w\left(v_{1}, \ldots, v_{n}\right)$.

Proposition $14.8 f_{1}, \ldots, f_{n}$ generate $\tilde{N}^{-}$freely, and so $\tilde{N}^{-}$is isomorphic to $F L\left(f_{1}, \ldots, f_{n}\right)$.

Proof. Define $\phi: \tilde{N}^{-} \rightarrow[T(V)]$ by $\phi(w)=\theta(w) \cdot 1$. Thus

$$
\phi\left(w\left(f_{1}, \ldots, f_{n}\right)\right)=w\left(v_{1}, \ldots, v_{n}\right)
$$

Then $\phi$ is a Lie algebra homomorphism, since

$$
\begin{aligned}
\phi\left[w\left(f_{1}, \ldots, f_{n}\right), w^{\prime}\left(f_{1}, \ldots, f_{n}\right)\right] & =\left[w\left(v_{1}, \ldots, v_{n}\right), w^{\prime}\left(v_{1}, \ldots, v_{n}\right)\right] \\
& =\left[\phi\left(w\left(f_{1}, \ldots, f_{n}\right)\right), \phi\left(w^{\prime}\left(f_{1}, \ldots, f_{n}\right)\right)\right] .
\end{aligned}
$$

Now $T(V)=F\left(v_{1}, \ldots, v_{n}\right)$, the free associative algebra on $v_{1}, \ldots, v_{n}$. Thus the free Lie algebra $F L\left(v_{1}, \ldots, v_{n}\right)$ lies in $[T(V)]$ and consists of all Lie words in $v_{1}, \ldots, v_{n}$. Thus $F L\left(v_{1}, \ldots, v_{n}\right)$ is the image of $\phi$. Hence the homomorphism

$$
\phi: \tilde{N}^{-} \rightarrow F L\left(v_{1}, \ldots, v_{n}\right)
$$

is surjective. But there is a Lie algebra homomorphism

$$
\phi^{\prime}: F L\left(v_{1}, \ldots, v_{n}\right) \rightarrow \tilde{N}^{-}
$$

with $\phi^{\prime}\left(v_{i}\right)=f_{i}$. Moreover we have $\phi \circ \phi^{\prime}=1$ on $F L\left(v_{1}, \ldots, v_{n}\right)$ and $\phi^{\prime} \circ \phi=1$ on $\tilde{N}^{-}$. Thus $\phi, \phi^{\prime}$ are inverse isomorphisms and $\tilde{N}^{-}$is isomorphic to $F L\left(f_{1}, \ldots, f_{n}\right)$.

Corollary $14.9 e_{1}, \ldots, e_{n}$ generate $\tilde{N}$ freely.
Proof. Apply the automorphism $\tilde{w}$ of Proposition 14.5. We have $\tilde{w}\left(\tilde{N}^{-}\right)=\tilde{N}$ and $\tilde{w}\left(f_{i}\right)=-e_{i}$. Thus the result follows from Proposition 14.8.

Proposition 14.10 $\tilde{L}(A)=\tilde{N}^{-} \oplus \tilde{H} \oplus \tilde{N}$, a direct sum of subspaces.

Proof. The proof is similar to that of Proposition 7.12. We show that $I=\tilde{N}^{-}+\tilde{H}+\tilde{N}$ is an ideal of $\tilde{L}(A)$. It is sufficient to show that

$$
\operatorname{ad} e_{i} \cdot I \subset I, \quad \text { ad } f_{i} \cdot I \subset I, \quad \operatorname{ad} \tilde{x} \cdot I \subset I
$$

Since the defining relations show that

$$
\begin{aligned}
& \operatorname{ad} e_{i} \cdot \tilde{H} \subset \tilde{N}, \quad \text { ad } e_{i} \cdot \tilde{N} \subset \tilde{N} \\
& \operatorname{ad} f_{i} \cdot \tilde{H} \subset \tilde{N}^{-}, \quad \operatorname{ad} f_{i} \cdot \tilde{N}^{-} \subset \tilde{N}^{-} \\
& \operatorname{ad} \tilde{x} \cdot \tilde{H}=O, \quad \operatorname{ad} \tilde{x} \cdot \tilde{N} \subset \tilde{N}, \quad \text { ad } \tilde{x} \cdot \tilde{N}^{-} \subset \tilde{N}^{-}
\end{aligned}
$$

it is sufficient to check that

$$
\begin{aligned}
& \operatorname{ad} f_{i} \cdot \tilde{N} \subset \tilde{H}+\tilde{N} \\
& \operatorname{ad} e_{i} \cdot \tilde{N}^{-} \subset \tilde{H}+\tilde{N}^{-}
\end{aligned}
$$

We have

$$
\operatorname{ad} f_{i} \cdot e_{j}=\delta_{i j} \tilde{h}_{i} \in \tilde{H}+\tilde{N}
$$

Suppose $w_{1}, w_{2} \in \tilde{N}$ satisfy

$$
\operatorname{ad} f_{i} \cdot w_{1} \in \tilde{H}+\tilde{N}, \quad \operatorname{ad} f_{i} \cdot w_{2} \in \tilde{H}+\tilde{N}
$$

Then

$$
\operatorname{ad} f_{i}\left[w_{1} w_{2}\right]=\left[\operatorname{ad} f_{i} \cdot w_{1}, w_{2}\right]+\left[w_{1}, \operatorname{ad} f_{i} \cdot w_{2}\right] \in \tilde{H}+\tilde{N}
$$

Thus ad $f_{i} \cdot \tilde{N} \subset \tilde{H}+\tilde{N}$.
The relation ad $e_{i} \cdot \tilde{N}^{-} \subset \tilde{H}+\tilde{N}^{-}$follows similarly. Thus $I$ is an ideal of $\tilde{L}(A)$ containing all the generators, and so $\tilde{L}(A)=\tilde{N}^{-}+\tilde{H}+\tilde{N}$.

In order to show the sum is direct we verify that if $w_{-} \in \tilde{N}^{-}, \tilde{x} \in \tilde{H}, w \in \tilde{N}$ satisfy

$$
w_{-}+\tilde{x}+w=0
$$

then we have $w_{-}=0, \tilde{x}=0, w=0$. Thus suppose $w_{-}+\tilde{x}+w=0$. Then $\theta_{\lambda}\left(w_{-}+\tilde{x}+w\right)$ is the zero endomorphism of $T(V)$. In particular $\theta_{\lambda}\left(w_{-}+\right.$ $\tilde{x}+w) \cdot 1=0$. Now $\theta_{\lambda}\left(w_{-}\right) \cdot 1=\phi\left(w_{-}\right), \theta_{\lambda}(\tilde{x}) \cdot 1=\lambda(x) 1$ and $\theta_{\lambda}(w) \cdot 1=0$. Hence

$$
\phi\left(w_{-}\right)+\lambda(x) 1=0
$$

Now $\phi\left(w_{-}\right) \in \bigoplus_{s \geq 1} T^{s}(V)$ and $\lambda(x) 1 \in T^{0}(V)$. It follows that $\phi\left(w_{-}\right)=0$ and $\lambda(x) 1=0$, that is $\lambda(x)=0$. Since this holds for all $\lambda \in H^{*}$ we have $x=0$. Hence $\tilde{x}=0$.

Now $\phi: \tilde{N}^{-} \rightarrow F L\left(v_{1}, \ldots, v_{n}\right)$ is an isomorphism, and so $\phi\left(w_{-}\right)=0$ implies $w_{-}=0$. Finally $w_{-}+\tilde{x}+w=0$ implies $w=0$. Thus

$$
\tilde{L}(A)=\tilde{N}^{-} \oplus \tilde{H} \oplus \tilde{N}
$$

Let $Q$ be the subgroup of $H^{*}$ given by $Q=\left\{\alpha=k_{1} \alpha_{1}+\cdots+k_{n} \alpha_{n} ; k_{1}, \ldots\right.$, $\left.k_{n} \in \mathbb{Z}\right\}$. Let $Q^{+}=\left\{\alpha \neq 0 \in Q ; k_{i} \geq 0\right.$ for all $\left.i\right\}$ and $Q^{-}=\left\{\alpha \neq 0 \in Q ; k_{i} \leq 0\right.$ for all $i\}$. For each $\alpha \in Q$ let

$$
\tilde{L}_{\alpha}=\{y \in \tilde{L}(A) ; \quad[\tilde{x} y]=\alpha(x) y \quad \text { for all } x \in H\}
$$

Proposition 14.11 (i) $\tilde{L}(A)=\oplus_{\alpha \in Q} \tilde{L}_{\alpha}$
(ii) $\operatorname{dim} \tilde{L}_{\alpha}$ is finite for all $\alpha \in Q$.
(iii) $\tilde{L}_{0}=\tilde{H}$.
(iv) If $\alpha \neq 0$ then $\tilde{L}_{\alpha}=0$ unless $\alpha \in Q^{+}$or $\alpha \in Q^{-}$.
(v) $\left[\tilde{L}_{\alpha} \tilde{L}_{\beta}\right] \subset \tilde{L}_{\alpha+\beta}$ for all $\alpha, \beta \in Q$.

Proof. To show $\tilde{L}(A)=\sum_{\alpha \in Q} \tilde{L}_{\alpha}$ it is sufficient to show $\tilde{H} \subset \sum_{\alpha \in Q} \tilde{L}_{\alpha}, \tilde{N} \subset$ $\sum_{\alpha \in Q} \tilde{L}_{\alpha}, \tilde{N}^{-} \subset \sum_{\alpha \in Q} \tilde{L}_{\alpha}$. It is clear that $\tilde{H} \subset \tilde{L}_{0}$. To show that $\tilde{N} \subset \sum_{\alpha \in Q^{+}} L_{\alpha}$ we observe that each Lie monomial $w$ in $e_{1}, \ldots, e_{n}$ satisfies $[\tilde{x} w]=\alpha(x) w$ for all $x \in H$ and some $\alpha \in Q^{+}$. For

$$
\left[\tilde{x} e_{i}\right]=\alpha_{i}(x) e_{i}
$$

and if

$$
\left[\tilde{x} w_{1}\right]=\beta(x) w_{1}, \quad\left[\tilde{x} w_{2}\right]=\gamma(x) w_{2}
$$

we have

$$
\left[\tilde{x}\left[w_{1} w_{2}\right]\right]=(\beta+\gamma)(x)\left[w_{1} w_{2}\right]
$$

This shows $\tilde{N} \subset \sum_{\alpha \in Q^{+}} \tilde{L}_{\alpha}$ and similarly we have $\tilde{N}^{-} \subset \sum_{\alpha \in Q^{-}} \tilde{L}_{\alpha}$. Thus $\tilde{L}(A)=\sum_{\alpha \in Q} \tilde{L}_{\alpha}$.

In order to show that the sum is direct we show that

$$
v_{1}+\cdots+v_{k}=0
$$

for $v_{i} \in \tilde{L}_{\beta_{i}}$ with $\beta_{1}, \ldots, \beta_{k}$ distinct implies each $v_{i}=0$. Suppose this is false. Choose the minimal value of $k$ for which it is false. Suppose $v_{1}+\cdots+v_{k}=0$ for this value of $k$ but that not each $v_{i}=0$. Then

$$
\left[\tilde{x}, v_{1}+\cdots+v_{k}\right]=0 \quad \text { for all } x \in H
$$

Thus

$$
\beta_{1}(x) v_{1}+\cdots+\beta_{k}(x) v_{k}=0
$$

We also have

$$
\beta_{k}(x) v_{1}+\cdots+\beta_{k}(x) v_{k}=0
$$

Hence

$$
\left(\beta_{1}(x)-\beta_{k}(x)\right) v_{1}+\cdots+\left(\beta_{k-1}(x)-\beta_{k}(x)\right) v_{k-1}=0 .
$$

By the minimality of $k$ we have

$$
\left(\beta_{i}(x)-\beta_{k}(x)\right) v_{i}=0 \quad \text { for } i=1, \ldots, k-1 .
$$

Since $\beta_{i} \neq \beta_{k}$ there exists $x \in H$ with $\beta_{i}(x) \neq \beta_{k}(x)$. Hence $v_{i}=0$ for $i=1, \ldots, k-1$. It follows that $v_{k}=0$. This contradicts our assumption. Hence

$$
\tilde{L}(A)=\bigoplus_{\alpha \in Q} \tilde{L}_{\alpha} .
$$

Since $\tilde{L}(A)=\tilde{N}^{-} \oplus \tilde{H} \oplus \tilde{N}$ by Proposition 14.10 and

$$
\tilde{N}^{-} \subset \sum_{\alpha \in Q^{-}} \tilde{L}_{\alpha}, \quad \tilde{H} \subset \tilde{L}_{0}, \quad \tilde{N} \subset \sum_{\alpha \in Q^{+}} \tilde{L}_{\alpha}
$$

it follows that

$$
\tilde{H}=\tilde{L}_{0}, \quad \tilde{N}=\sum_{\alpha \in Q^{+}} \tilde{L}_{\alpha}, \quad \tilde{N}^{-}=\sum_{\alpha \in Q^{-}} \tilde{L}_{\alpha} .
$$

Also we have $\tilde{L}_{\alpha}=O$ if $\alpha \neq 0, \alpha \notin Q^{+}, \alpha \notin Q^{-}$. The Jacobi identity shows that $\left[\tilde{L}_{\alpha} \tilde{L}_{\beta}\right] \subset \tilde{L}_{\alpha+\beta}$ for all $\alpha, \beta \in Q$.
Finally we show $\operatorname{dim} \tilde{L}_{\alpha}$ is finite. We have $\operatorname{dim} \tilde{L}_{0}=2 n-l$. So let $\alpha \in Q^{+}$. Then $\tilde{L}_{\alpha} \subset \tilde{N}$. Now $\tilde{N}$ is spanned by Lie monomials in $e_{1}, \ldots, e_{n}$ and each Lie monomial lies in some $\tilde{L}_{\alpha}$. Let $\alpha=k_{1} \alpha_{1}+\cdots+k_{n} \alpha_{n}$ with $k_{i} \in \mathbb{Z}$ and $k_{i} \geq 0$. A Lie monomial lies in $\tilde{L}_{\alpha}$ if and only if $e_{i}$ appears $k_{i}$ times in it for each $i$. But there are only finitely many Lie monomials in which $e_{i}$ appears $k_{i}$ times for each $i$. Thus $\operatorname{dim} \tilde{L}_{\alpha}$ is finite. A similar argument proves this when $\alpha \in Q^{-}$. We note in particular that

$$
\begin{aligned}
& \operatorname{dim} \tilde{L}_{\alpha_{i}}=1, \quad \operatorname{dim} \tilde{L}_{-\alpha_{i}}=1 \\
& \operatorname{dim} \tilde{L}_{k \alpha_{i}}=0, \quad \operatorname{dim} \tilde{L}_{-k \alpha_{i}}=0 \quad \text { if } k>1 .
\end{aligned}
$$

The following lemma will be needed in the proof of the next proposition.
Lemma 14.12 Let $H$ be a finite dimensional abelian Lie algebra and $V$ be an H -module such that

$$
V=\bigoplus_{\lambda \in H^{*}} V_{\lambda}
$$

where $V_{\lambda}=\{v \in V ; x v=\lambda(x) v$ for all $x \in H\}$. Let $U$ be a submodule of $V$. Then

$$
U=\bigoplus_{\lambda \in H^{*}}\left(U \cap V_{\lambda}\right) .
$$

Proof. Let $u \in U$. Then $u=u_{1}+\cdots+u_{m}$ where $u_{i} \in V_{\lambda_{i}}$ and $\lambda_{1}, \ldots, \lambda_{m}$ are distinct elements of $H^{*}$. Let

$$
H_{i j}=\left\{x \in H ; \quad \lambda_{i}(x)=\lambda_{j}(x)\right\} \quad \text { for } i \neq j
$$

$H_{i j}$ is a subspace of $H$ of codimension 1 . Now $H \neq \bigcup_{i \neq j} H_{i j}$ since a finite dimensional vector space over $\mathbb{C}$ cannot be the union of finitely many proper subspaces. So we can find $x \in H$ with $\lambda_{1}(x), \ldots, \lambda_{m}(x)$ all distinct.

Let $\theta(x): V \rightarrow V$ be the linear map given by $\theta(x) v=x v$. Then we have

$$
\begin{aligned}
& u=u_{1}+\cdots+u_{m} \\
& \theta(x) u=\lambda_{1}(x) u_{1}+\cdots+\lambda_{m}(x) u_{m} \\
& \theta(x)^{2} u=\lambda_{1}(x)^{2} u_{1}+\cdots+\lambda_{m}(x)^{2} u_{m} \\
& \vdots \\
& \theta(x)^{m-1} u=\lambda_{1}(x)^{m-1} u_{1}+\cdots+\lambda_{m}(x)^{m-1} u_{m} .
\end{aligned}
$$

We have here $m$ equations in $u_{1}, \ldots, u_{m}$ whose coefficients have non-zero determinant. Thus $u_{1}, \ldots, u_{m}$ may be expressed as linear combinations of $u, \theta(x) u, \theta(x)^{2} u, \ldots, \theta(x)^{m-1} u$. These vectors all lie in $U$. Thus $u_{i} \in U \cap V_{\lambda_{i}}$. Thus we have shown that $U=\sum_{\lambda \in H^{*}}\left(U \cap V_{\lambda}\right)$ and the sum is direct because $\sum_{\lambda \in H^{*}} V_{\lambda}$ is a direct sum.

Proposition 14.13 The algebra $\tilde{L}(A)$ contains a unique ideal I maximal with respect to $I \cap \tilde{H}=O$.

Proof. Let $J$ be any ideal of $\tilde{L}(A)$ with $J \cap \tilde{H}=O$. We have

$$
\tilde{L}(A)=\bigoplus_{\alpha \in H^{*}} \tilde{L}_{\alpha}
$$

by Proposition 14.11, and we consider $\tilde{L}(A)$ as an $\tilde{H}$-module. By Lemma 14.12 we have

$$
J=\bigoplus_{\alpha \in H^{*}}\left(\tilde{L}_{\alpha} \cap J\right)
$$

Now each $\tilde{L}_{\alpha}$ with $\alpha \neq 0$ lies in $\tilde{N}$ or in $\tilde{N}^{-}$. Thus

$$
J=\left(\tilde{N}^{-} \cap J\right) \oplus(\tilde{N} \cap J)
$$

In particular $J \subset \tilde{N}^{-} \oplus \tilde{N}$.
Now consider the ideal $I$ of $\tilde{L}(A)$ generated by all ideals $J$ with $J \cap \tilde{H}=O$. All such ideals $J$ lie in $\tilde{N}^{-} \oplus \tilde{N}$, thus $I$ lies in $\tilde{N}^{-} \oplus \tilde{N}$. Hence $I \cap \tilde{H}=O$. Thus $I$ is the unique ideal of $\tilde{L}(A)$ maximal with respect to $I \cap \tilde{H}=O$.

### 14.3 The Kac-Moody algebra $L(A)$

We now suppose that $A$ is a GCM. Let $\tilde{L}(A)$ be the Lie algebra associated with $A$ defined in Section 14.2 and $I$ be the unique maximal ideal of $\tilde{L}(A)$ with $I \cap \tilde{H}=O$. Let $L(A)$ be defined by

$$
L(A)=\tilde{L}(A) / I
$$

The Lie algebra $L(A)$ is called the Kac-Moody algebra with GCM $A$. We have a natural homomorphism $\theta: \tilde{L}(A) \rightarrow L(A)$. We define $N=\theta(\tilde{N})$ and $N^{-}=\theta\left(\tilde{N}^{-}\right)$.

Proposition 14.14 $L(A)=N^{-} \oplus \theta(\tilde{H}) \oplus N$. Moreover the map $\theta: \tilde{H} \rightarrow \theta(\tilde{H})$ is an isomorphism.

Proof. We know from the proof of Proposition 14.13 that

$$
I=\left(\tilde{N}^{-} \cap I\right) \oplus(\tilde{N} \cap I)
$$

Since $\tilde{L}(A)=\tilde{N}^{-} \oplus \tilde{H} \oplus \tilde{N}$ it follows that

$$
L(A)=N^{-} \oplus \theta(\tilde{H}) \oplus N
$$

and that $\theta: \tilde{H} \rightarrow \theta(\tilde{H})$ is an isomorphism.
We recall from Corollary 14.7 that there is a natural isomorphism $H \rightarrow \tilde{H}$. Combining this with $\theta$ we obtain an isomorphism $H \rightarrow \theta(\tilde{H})$. We shall subsequently use this isomorphism to identify $\theta(\tilde{H})$ with $H$, and we shall write

$$
L(A)=N^{-} \oplus H \oplus N
$$

In order to show that a given Lie algebra is isomorphic to $L(A)$ the following result is often useful.

Proposition 14.15 Suppose we are given an $n \times n G C M A=\left(A_{i j}\right)$. Let L be a Lie algebra over $\mathbb{C}$ and $H$ be a finite dimensional abelian subalgebra of $L$
with $\operatorname{dim} H=2 n-\operatorname{rank} A$. Suppose $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ is a linearly independent subset of $H^{*}$ and $\Pi^{v}=\left\{h_{1}, \ldots, h_{n}\right\}$ a linearly independent subset of $H$ satisfying $\alpha_{j}\left(h_{i}\right)=A_{i j}$.

Suppose also that $e_{1}, \ldots, e_{n}, f_{1}, \ldots, f_{n}$ are elements of $L$ satisfying

$$
\begin{aligned}
& {\left[e_{i} f_{i}\right]=h_{i}} \\
& {\left[e_{i} f_{j}\right]=0 \quad \text { if } i \neq j} \\
& {\left[x e_{i}\right]=\alpha_{i}(x) e_{i} \quad \text { for } x \in H} \\
& {\left[x f_{i}\right]=-\alpha_{i}(x) f_{i} \quad \text { for } x \in H}
\end{aligned}
$$

Suppose that $e_{1}, \ldots, e_{n}, f_{1}, \ldots, f_{n}$ and $H$ generate $L$ and that $L$ has no non-zero ideal $J$ with $J \cap H=O$. Then $L$ is isomorphic to the Kac-Moody algebra $L(A)$.

Proof. The elements $e_{1}, \ldots, e_{n}, f_{1}, \ldots, f_{n}$ and $x \in H$ generate $L$ and satisfy all the defining relations of $\tilde{L}(A)$ given in Section 14.2. Thus there is a surjective Lie algebra homomorphism $\theta: \tilde{L}(A) \rightarrow L$ and $L$ is isomorphic to $\tilde{L}(A) / \operatorname{ker} \theta$. The restriction map $\theta: \tilde{H} \rightarrow H$ is an isomorphism by Corollary 14.7, thus $\operatorname{ker} \theta \cap \tilde{H}=O$. It follows that $\operatorname{ker} \theta \subset I$, the largest ideal of $\tilde{L}(A)$ with $I \cap \tilde{H}=O$. In fact we have $\operatorname{ker} \theta=I$ since $L$ has no non-zero ideal $J$ with $J \cap H=O$. Hence

$$
L \cong \tilde{L}(A) / I=L(A)
$$

Corollary 14.16 If $A$ is a Cartan matrix then $L(A)$ is the finite dimensional semisimple Lie algebra with Cartan matrix A.

Proof. In this case we have rank $A=n$, so $\operatorname{dim} H=n$. The finite dimensional semisimple Lie algebra satisfies all the hypotheses of Proposition 14.15, so is isomorphic to the Kac-Moody algebra $L(A)$.

This result shows that the theory of Kac-Moody algebras is an extension of the theory of finite dimensional semisimple Lie algebras, which we have already described.

We shall now describe some further basic properties of the Kac-Moody algebra $L(A)$. We shall denote the images of $e_{i}, h_{i}, f_{i} \in \tilde{L}(A)$ under the natural homomorphism $\tilde{L}(A) \rightarrow L(A)$ by $e_{i}, h_{i}, f_{i} \in L(A)$. This should not lead to confusion as we shall subsequently be concentrating on $L(A)$ rather than $\tilde{L}(A)$.

Proposition 14.17 There is an automorphism $\omega$ of $L(A)$ satisfying $\omega^{2}=1$ determined by

$$
\begin{array}{ll}
\omega\left(e_{i}\right)=-f_{i}, & \omega\left(f_{i}\right)=-e_{i} \\
\omega(x)=-x & \text { for all } x \in H
\end{array}
$$

Proof. By Proposition $14.5 \tilde{L}(A)$ has an automorphism $\tilde{\omega}$ with $\tilde{\omega}^{2}=1$. Thus $\tilde{\omega}(I)$ is the unique maximal ideal with

$$
\tilde{\omega}(I) \cap \tilde{\omega}(\tilde{H})=O .
$$

But $\tilde{\omega}(\tilde{H})=\tilde{H}$ so $\tilde{\omega}(I)=I$. Thus $\tilde{\omega}$ induces an automorphism $\omega$ of $\tilde{L}(A) / I=$ $L(A)$ satisfying the stated conditions.

There is also an analogue of Proposition 14.11. For each $\alpha \in Q$ define $L_{\alpha}$ by

$$
L_{\alpha}=\{y \in L(A) ;[x y]=\alpha(x) y \quad \text { for all } x \in H\}
$$

Proposition 14.18 (i) $L(A)=\bigoplus_{\alpha \in Q} L_{\alpha}$
(ii) $\operatorname{dim} L_{\alpha}$ is finite for all $\alpha \in Q$.
(iii) $L_{0}=H$
(iv) If $\alpha \neq 0$ then $L_{\alpha}=O$ unless $\alpha \in Q^{+}$or $\alpha \in Q^{-}$.
(v) $\left[L_{\alpha} L_{\beta}\right] \subset L_{\alpha+\beta}$ for all $\alpha, \beta \in Q$.

Proof. Let $\theta: \tilde{L}(A) \rightarrow L(A)=\tilde{L}(A) / I$ be the natural homomorphism. We have

$$
\tilde{L}(A)=\bigoplus_{\alpha \in Q} \tilde{L}_{\alpha} \quad \text { by Proposition } 14.11
$$

Also

$$
I=\bigoplus_{\alpha \in Q}\left(I \cap \tilde{L}_{\alpha}\right) \quad \text { by Lemma } 14.12
$$

It follows that

$$
L(A)=\bigoplus_{\alpha \in Q} \theta\left(\tilde{L}_{\alpha}\right)
$$

Now we clearly have $\theta\left(\tilde{L}_{\alpha}\right) \subset L_{\alpha}$, thus $L(A)=\sum_{\alpha \in Q} L_{\alpha}$. This sum is direct, just as in the proof of Proposition 14.11. It follows that $L(A)=\bigoplus_{\alpha \in Q} L_{\alpha}$ and that $L_{\alpha}=\theta\left(\tilde{L}_{\alpha}\right)$. Now

$$
L(A)=N^{-} \oplus H \oplus N \quad \text { by Proposition } 14.14
$$

and $N^{-} \subset \sum_{\alpha \in Q^{-}} L_{\alpha}, \quad H \subset L_{0}, \quad N \subset \sum_{\alpha \in Q^{+}} L_{\alpha}$, hence we have

$$
N^{-}=\bigoplus_{\alpha \in Q^{-}} L_{\alpha}, \quad H=L_{0}, \quad N=\bigoplus_{\alpha \in Q^{+}} L_{\alpha}
$$

$\operatorname{dim} L_{\alpha}$ is finite because $L_{\alpha}=\theta\left(\tilde{L}_{\alpha}\right)$ and $\operatorname{dim} \tilde{L}_{\alpha}$ is finite. Finally $\left[L_{\alpha} L_{\beta}\right] \subset$ $L_{\alpha+\beta}$ follows from the Jacobi identity.

Definitions $H$ will be called a Cartan subalgebra of $L(A)$. This fits in with our previous terminology when $A$ was a Cartan matrix. An element $\alpha \in H^{*}$ is called a root of $L(A)$ if $\alpha \neq 0$ and $L_{\alpha} \neq O$. Every root lies in $Q^{+}$or $Q^{-}$. The roots in $Q^{+}$are called positive roots and those in $Q^{-}$negative roots. If $\alpha$ is a root then $L_{\alpha}$ is called the root space of $\alpha$. The dimension of $L_{\alpha}$ is called the multiplicity of $\alpha$. When A is a Cartan matrix we recall that all roots have multiplicity 1. However, we shall see that this is not always the case when $A$ is a GCM.

Proposition 14.19 (i) $\operatorname{dim} L_{\alpha_{i}}=1$ and $\operatorname{dim} L_{-\alpha_{i}}=1$.
(ii) If $k>1$ then $\operatorname{dim} L_{k \alpha_{i}}=0, \operatorname{dim} L_{-k \alpha_{i}}=0$.

Proof. Since $L_{\alpha_{i}}=\theta\left(\tilde{L}_{\alpha_{i}}\right)$ and $\operatorname{dim} \tilde{L}_{\alpha_{i}}=1$ we have $\operatorname{dim} L_{\alpha_{i}} \leq 1$. If $\operatorname{dim} L_{\alpha_{i}}=0$ we would have $e_{i} \in I=\operatorname{ker} \theta$. This would imply $\left[e_{i} f_{i}\right]=\tilde{h}_{i} \in I$, contrary to $I \cap \tilde{H}=O$. Thus $\operatorname{dim} L_{\alpha_{i}}=1$. A similar argument gives $\operatorname{dim} L_{-\alpha_{i}}=1$.

Since $\tilde{L}_{k \alpha_{i}}=O$ and $\tilde{L}_{-k \alpha_{i}}=O$ for $k>1$ it follows that $L_{k \alpha_{i}}=O$ and $L_{-k \alpha_{i}}=O$.
$\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n}$ are called the fundamental roots of $L(A)$, again in agreement with the earlier terminology when $A$ is a Cartan matrix.

Remark 14.20 For a general $n \times n$ matrix $A$ over $\mathbb{C}$ we constructed a minimal realisation $\left(H, \Pi, \Pi^{\mathrm{v}}\right)$ where $H$ is a vector space over $\mathbb{C}$ of dimension $2 n-\operatorname{rank} A, \Pi^{\mathrm{v}}=\left\{h_{1}, \ldots, h_{n}\right\}$ is a linearly independent subset of $H$ and $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ is a linearly independent subset of $H^{*}$ such that $\alpha_{j}\left(h_{i}\right)=A_{i j}$.

In the case when $A$ is a GCM the matrix $A$ is real and so we can find a real vector space $H_{\mathbb{R}}$, of dimension $2 n$ - rank $A$ over $\mathbb{R}$, contained in $H$ such that $h_{1}, \ldots, h_{n}$ lie in $H_{\mathbb{R}}$ and are linearly independent and $\alpha_{1}, \ldots, \alpha_{n}$, when restricted to $H_{\mathbb{R}}^{*}$, remain linearly independent. In the construction of $H$, described in Proposition 14.2 as the vector space of all $(2 n-l)$-tuples over $\mathbb{C}$, we define $H_{\mathbb{R}}$ as the subset of all $(2 n-l)$-tuples over $\mathbb{R}$. The triple ( $H_{\mathbb{R}}, \Pi, \Pi^{\mathrm{v}}$ ) with $\Pi^{\mathrm{v}} \subset H_{\mathbb{R}}$ and $\Pi \subset H_{\mathbb{R}}^{*}$ is called a real minimal realisation of $A$.

We denote by $L(A)^{\prime}$ the subalgebra of $L(A)$ generated by $e_{1}, \ldots, e_{n}$, $f_{1}, \ldots, f_{n}$.

Proposition 14.21 (i) $L_{\alpha}$ lies in $L(A)^{\prime}$ for each root $\alpha$ of $L(A)$.
(ii) $L(A)^{\prime}=\left(H \cap L(A)^{\prime}\right) \oplus \sum_{\alpha \neq 0} L_{\alpha}$.
(iii) $L(A)^{\prime}=[L(A) L(A)]$.

Proof. We know from Proposition 14.18 that $L_{\alpha} \neq O$ implies $\alpha \in Q^{+}$or $\alpha \in$ $Q^{-}$. If $\alpha \in Q^{+}$then $L_{\alpha} \subset N$ and if $\alpha \in Q^{-}$then $L_{\alpha} \subset N^{-}$. Since $N$ is the subalgebra generated by $e_{1}, \ldots, e_{n}$ and $N^{-}$is the subalgebra generated by $f_{1}, \ldots, f_{n}$ we have $L_{\alpha} \subset L(A)^{\prime}$ for each $\alpha$.

Since $L(A)=H \oplus \sum_{\alpha \neq 0} L_{\alpha}$ and $L_{\alpha} \subset L(A)^{\prime}$ we have

$$
L(A)^{\prime}=\left(H \cap L(A)^{\prime}\right) \oplus \sum_{\alpha \neq 0} L_{\alpha}
$$

It follows that $L(A)=L(A)^{\prime}+H$. We also have $\left[H, L(A)^{\prime}\right] \subset L(A)^{\prime}$ and so $L(A)^{\prime}$ is an ideal of $L(A)$. We have

$$
L(A) / L(A)^{\prime} \cong H / H \cap L(A)^{\prime}
$$

and so $L(A) / L(A)^{\prime}$ is abelian. Hence $[L(A) L(A)] \subset L(A)^{\prime}$. On the other hand we have $\left[e_{i} f_{i}\right]=h_{i},\left[h_{i} e_{i}\right]=2 e_{i},\left[h_{i} f_{i}\right]=-2 f_{i}$ and so $e_{i}, f_{i} \in[L(A) L(A)]$. Thus $L(A)^{\prime} \subset[L(A) L(A)]$ and we have equality.

## 15

## The classification of generalised Cartan matrices

The structure of the Kac-Moody algebra $L(A)$ depends crucially on the GCM $A$. In the present chapter we shall discuss various possible types of GCM $A$ which can occur.

### 15.1 A trichotomy for indecomposable GCMs

Two GCMs $A, A^{\prime}$ are called equivalent if they have the same degree $n$ and there is a permutation $\sigma$ of $1, \ldots, n$ such that

$$
A_{i j}^{\prime}=A_{\sigma(i) \sigma(j)} \quad \text { for all } i, j .
$$

A GCM $A$ is called indecomposable if it is not equivalent to a diagonal sum

$$
\left(\begin{array}{ll}
A_{1} & O \\
O & A_{2}
\end{array}\right)
$$

of smaller GCMs $A_{1}, A_{2}$. If $A$ is a GCM so is its transpose $A^{\mathrm{t}}$. Moreover $A$ is indecomposable if and only if $A^{\mathrm{t}}$ is indecomposable.

We shall now define three particular types of GCM. Let $v=\left(v_{1}, \ldots, v_{n}\right)$ be a vector in $\mathbb{R}^{n}$. We write $v \geq 0$ if $v_{i} \geq 0$ for each $i$, and $v>0$ if $v_{i}>0$ for each $i$.

Definitions A GCM A has finite type if
(i) $\operatorname{det} A \neq 0$
(ii) there exists $u>0$ with $A u>0$
(iii) $A u \geq 0$ implies $u>0$ or $u=0$.

## The GCM A has affine type if

(i) corank $A=1 \quad$ (i.e. rank $A=n-1)$
(ii) there exists $u>0$ such that $A u=0$
(iii) $A u \geq 0$ implies $A u=0$.

The GCM A has indefinite type if
(i) there exists $u>0$ such that $A u<0$
(ii) $A u \geq 0$ and $u \geq 0$ imply $u=0$.

All vectors $u$ in these definitions are assumed to lie in $\mathbb{R}^{n}$, and are column vectors.

We aim to prove the following theorem.

Theorem 15.1 Let A be an indecomposable GCM. Then exactly one of the following three possibilities holds:
(a) A has finite type
(b) A has affine type
(c) A has indefinite type.

Moreover the type of $A^{\mathrm{t}}$ is the same as the type of $A$.
This section will be devoted to the proof of Theorem 15.1, which gives a trichotomy on the set of indecomposable GCMs.

We begin with a lemma on inequalities.

Lemma 15.2 Let $v^{i}=\left(v_{i 1}, \ldots, v_{i n}\right) \in \mathbb{R}^{n}$ for $i=1, \ldots$, . Then there exist $x_{1}, \ldots, x_{n} \in \mathbb{R}$ with $\sum_{j=1}^{n} v_{i j} x_{j}>0$ for $i=1, \ldots, m$ if and only if $\lambda_{1} v^{1}+\cdots+$ $\lambda_{m} v^{m}=0, \lambda_{i} \geq 0$ implies $\lambda_{i}=0$ for $i=1, \ldots, m$.

Proof. Suppose there exists a column vector $x=\left(x_{1}, \ldots, x_{n}\right)^{\mathrm{t}}$ satisfying $v^{i} x>$ 0 for all $i$. Suppose $\lambda_{1} v^{1}+\cdots+\lambda_{m} v^{m}=0$ with all $\lambda_{i} \geq 0$. Then $\lambda_{1} v^{1} x+\cdots+$ $\lambda_{m} v^{m} x=0$. But $v^{i} x>0$ and $\lambda_{i} \geq 0$, thus we have $\lambda_{i}=0$ for all $i$.

Conversely suppose $\lambda_{1} v^{1}+\cdots+\lambda_{m} v^{m}=0, \lambda_{i} \geq 0$ implies $\lambda_{i}=0$ for all $i$. Let

$$
S=\left\{\sum_{i=1}^{m} \lambda_{i} v^{i} ; \quad \lambda_{i} \geq 0, \quad \sum_{i=1}^{m} \lambda_{i}=1\right\}
$$

Define $f: S \rightarrow \mathbb{R}$ by $f(y)=\|y\|$ where $\|y\|=\sqrt{y_{1}^{2}+\cdots+y_{n}^{2}}$. Then $S$ is a compact subset of $\mathbb{R}^{n}$ and $f$ is a continuous function from $S$ to $\mathbb{R}$. Thus $f(S)$ is a compact subset of $\mathbb{R}$. Hence there exists $x \in S$ with $\|x\| \leq\left\|x^{\prime}\right\|$ for all $x^{\prime} \in S$. Clearly $x \neq 0$ since the zero vector does not lie in $S$. We shall show
$v^{i} x>0$ for all $i$ as required. In fact we shall show that $(y, x)>0$ for all $y \in S$, where $(y, x)=\sum y_{i} x_{i}$. This implies the required result since each $v^{i}$ lies in $S$.

Now $S$ is a convex subset of $\mathbb{R}^{n}$. We assume $y \neq x$, then $t y+(1-t) x \in S$ for all $t$ with $0 \leq t \leq 1$. By the choice of $x$ we have

$$
(t y+(1-t) x, \quad t y+(1-t) x) \geq(x, x)
$$

that is

$$
t(y-x, y-x)+2(y-x, x) \geq 0
$$

for $0<t \leq 1$. This implies $(y-x, x) \geq 0$, that is $(y, x) \geq(x, x)>0$.
We make use of this lemma in the following proposition.

Proposition 15.3 Let $M$ be an $m \times n$ matrix over $\mathbb{R}$. Suppose

$$
u \geq 0 \text { and } M^{\mathrm{t}} u \geq 0 \text { imply } u=0
$$

Then there exists $v>0$ with $M v<0$.
Proof. Let $M=\left(m_{i j}\right)$ and consider the following system of inequalities:

$$
\begin{aligned}
& -\sum_{j=1}^{n} m_{i j} x_{j}>0 \quad i=1, \ldots, m \\
& x_{j}>0 \quad j=1, \ldots, n
\end{aligned}
$$

We shall use Lemma 15.2 to show that these inequalities have a solution. Thus we consider an equation of form

$$
\sum_{i=1}^{m} \lambda_{i}\left(-m_{i 1}, \ldots,-m_{i n}\right)+\sum_{j=1}^{n} \mu_{j}\left(0, \ldots, 1_{j}, \ldots, 0\right)=0
$$

with $\lambda_{i} \geq 0, \mu_{j} \geq 0$ for all $i, j$. Then

$$
\sum_{i=1}^{m} \lambda_{i} m_{i j}=\mu_{j}
$$

Let $u=\left(\lambda_{1}, \ldots, \lambda_{m}\right)^{\mathrm{t}}$. Then $M^{\mathrm{t}} u=\left(\mu_{1}, \ldots, \mu_{n}\right)^{\mathrm{t}}$. Thus we have $u \geq 0$ and $M^{\mathrm{t}} u \geq 0$. This implies that $u=0$. We also have $M^{\mathrm{t}} u=0$. Thus $\lambda_{i}=0$ and $\mu_{j}=0$ for all $i, j$. Hence Lemma 15.2 shows that the above inequalities have a solution. Thus there exists $v>0$ with $M v<0$.

We now consider our three classes of GCM $A$. Let

$$
\begin{aligned}
S_{\mathrm{F}} & =\{A ; A \text { has finite type }\} \\
S_{\mathrm{A}} & =\{A ; A \text { has affine type }\} \\
S_{\mathrm{I}} & =\{A ; A \text { has indefinite type }\} .
\end{aligned}
$$

It is easy to see that no GCM can lie in more than one of these classes.

Lemma 15.4 $S_{\mathrm{F}} \cap S_{\mathrm{A}}=\phi, \quad S_{\mathrm{F}} \cap S_{\mathrm{I}}=\phi, \quad S_{\mathrm{A}} \cap S_{\mathrm{I}}=\phi$.

Proof. If $A \in S_{\mathrm{F}} \cap S_{\mathrm{A}}$ then $\operatorname{det} A \neq 0$ and corank $A=1$, a contradiction.
If $A \in S_{\mathrm{F}} \cap S_{\text {I }}$ there exists $u>0$ with $A u>0$. But $A u \geq 0$ and $u \geq 0$ imply $u=0$, a contradiction.

If $A \in S_{\mathrm{A}} \cap S_{\text {I }}$ there exists $u>0$ with $A u=0$. But $A u \geq 0$ and $u \geq 0$ imply $u=0$, a contradiction.

We must therefore show that each indecomposable GCM lies in one of the three classes.

Lemma 15.5 Let $A$ be an indecomposable GCM. Then $u \geq 0$ and $A u \geq 0$ imply that $u>0$ or $u=0$.

Proof. Suppose $u \neq 0$ and $u \ngtr 0$. Then we can reorder $1, \ldots, n$ so that $u_{i}=0$ for $i=1, \ldots, s$ and $u_{i}>0$ for $i=s+1, \ldots, n$. Let

$$
\begin{gathered}
A=\left(\begin{array}{cc}
P & Q \\
R & S
\end{array}\right)_{n-s}^{s} \\
s n-s
\end{gathered}
$$

Now all entries of the block $Q$ are $\leq 0$ and if $Q$ has an entry $<0$ then $A u$ has a negative coefficient, which is impossible. Thus $Q=0$. This implies $R=O$ by the definition of a GCM, thus $A$ is decomposable, a contradiction.

Now let $A$ be an indecomposable GCM and define $K_{A}$ by

$$
K_{A}=\{u ; A u \geq 0\}
$$

$K_{A}$ is a convex cone. We consider its intersection with the convex cone $\{u ; u \geq 0\}$. We shall distinguish between two cases:

$$
\begin{aligned}
& \{u ; u \geq 0, A u \geq 0\} \neq\{0\} \\
& \{u ; u \geq 0, A u \geq 0\}=\{0\} .
\end{aligned}
$$

The first of these cases splits into two subcases, as is shown by the next lemma.

Lemma 15.6 Suppose $\{u ; u \geq 0, A u \geq 0\} \neq\{0\}$. Then just one of the following cases occurs:

$$
\begin{aligned}
& K_{A} \subset\{u ; u>0\} \cup\{0\} \\
& K_{A}=\{u ; A u=0\} \quad \text { and } \quad K_{A} \quad \text { is a 1-dimensional subspace of } \mathbb{R}^{n} .
\end{aligned}
$$

Proof. We know there exists $u \neq 0$ with $u \geq 0$ and $A u \geq 0$. By Lemma 15.5 this implies that $u>0$. Suppose the first case does not hold. Then there exists $v \neq 0$ with $A v \geq 0$ such that some coordinate of $v$ is $\leq 0$. If $v \geq 0$ then $v>0$ by Lemma 15.5 , thus some coordinate of $v$ is $<0$.

We have $A u \geq 0$ and $A v \geq 0$, hence $A(t u+(1-t) v) \geq 0$ for $0 \leq t \leq 1$. Since all coordinates of $u$ are positive and some coordinate of $v$ is negative there exists $t$ with $0<t<1$ such that $t u+(1-t) v \geq 0$ and some coordinate of $t u+(1-t) v$ is 0 . But then $t u+(1-t) v=0$ by a further use of Lemma 15.5. Thus $v$ is a scalar multiple of $u$. We also have

$$
0=A(t u+(1-t) v)=t A u+(1-t) A v
$$

Since $A u \geq 0, A v \geq 0$ this implies that $A u=0, A v=0$.
Now let $w \in K_{A}$. Then $A w \geq 0$. Either $w \geq 0$ or some coordinate of $w$ is negative. If $w \geq 0$ then $w>0$ or $w=0$ by Lemma 15.5. Suppose $w>0$. Then by the above argument with $u$ replaced by $w, v$ is a scalar multiple of $w$. Hence $w$ is a scalar multiple of $u$. Now suppose some coordinate of $w$ is negative. Then by the above argument with $v$ replaced by $w, w$ is a scalar multiple of $u$. Thus in all cases $w$ is a scalar multiple of $u$. Hence $K_{A}$ is the 1-dimensional subspace $\mathbb{R} u$. Thus we have shown that $K_{A}$ is a 1-dimensional subspace of $\mathbb{R}^{n}$. We have also shown that if $w \in K_{A}$ then $A w=0$. Thus $K_{A}=\{w ; A w=0\}$.

Thus if the first case does not hold the second case must hold. We note finally that the two cases cannot hold together since in the first case $K_{A}$ cannot contain a 1 -dimensional subspace of $\mathbb{R}^{n}$.

We now identify the first case in Lemma 15.6 with the case of matrices of finite type.

Proposition 15.7 Let $A$ be an indecomposable GCM. Then the following conditions are equivalent:

A has finite type

$$
\{u ; u \geq 0 \text { and } A u \geq 0\} \neq\{0\} \quad \text { and } \quad K_{A} \subset\{u ; u>0\} \cup\{0\}
$$

Proof. Suppose $A$ has finite type. Then there exists $u>0$ with $A u>0$. Hence $\{u ; u \geq 0$ and $A u \geq 0\} \neq\{0\}$. Also $\operatorname{det} A \neq 0$. Thus $\{u ; A u=0\}$ is not a 1 -dimensional subspace of $\mathbb{R}^{n}$. Hence $K_{A} \subset\{u ; u>0\} \cup\{0\}$ by Lemma 15.6.

Conversely suppose $\{u ; u \geq 0$ and $A u \geq 0\} \neq\{0\}$ and $K_{A} \subset\{u ; u>0\} \cup$ $\{0\}$. Then there cannot exist $u \neq 0$ with $A u=0$. For this would give a 1 -dimensional subspace contained in $K_{A}$. Thus $\operatorname{det} A \neq 0$. Now there exists $u \neq 0$ with $u \geq 0$ and $A u \geq 0$. By Lemma 15.5 we have $u>0$. If $A u>0, A$ has finite type. So suppose to the contrary that some coordinates of $A u$ are zero and some are non-zero. We choose the numbering $1, \ldots, n$ so that the first $s$ components of $A u$ are 0 and the last $n-s$ are positive. Let

$$
\begin{gathered}
A=\left(\begin{array}{ll}
P & Q \\
R & S
\end{array}\right)_{n-s}^{s} \\
s \\
n-s
\end{gathered}
$$

Now the block $Q$ satisfies $Q \neq O$ since $A$ is indecomposable. We choose the numbering so that the first row of $Q$ is not the zero vector. We have

$$
A u=\left(\begin{array}{ll}
P & Q \\
R & S
\end{array}\right)\binom{u^{1}}{u^{2}}=\binom{P u^{1}+Q u^{2}}{R u^{1}+S u^{2}}
$$

Hence $P u^{1}+Q u^{2}=0$ and $R u^{1}+S u^{2}>0$. We also have $u^{1}>0, u^{2}>0$. Thus $Q u^{2} \leq 0$ and the first coordinate of $Q u^{2}$ is $<0$. Hence $P u^{1} \geq 0$ and the first coordinate of $P u^{1}$ is $>0$. Since $R u^{1}+S u^{2}>0$ we can choose $\varepsilon>0$ such that $R(1+\varepsilon) u^{1}+S u^{2}>0$.

We now consider, instead of our original vector $u=\binom{u^{1}}{u^{2}}$, the vector $\binom{(1+\varepsilon) u^{1}}{u^{2}}$. We have

$$
\begin{aligned}
& \binom{(1+\varepsilon) u^{1}}{u^{2}}>0 \\
& A\binom{(1+\varepsilon) u^{1}}{u^{2}}=\binom{P u^{1}+Q u^{2}+\varepsilon P u^{1}}{R u^{1}+S u^{2}+\varepsilon R u^{1}}=\binom{\varepsilon P u^{1}}{R(1+\varepsilon) u^{1}+S u^{2}}
\end{aligned}
$$

The first coordinate and the last $n-s$ coordinates of this vector are positive and the remaining coordinates are $\geq 0$. Thus

$$
A\binom{(1+\varepsilon) u^{1}}{u^{2}} \geq 0
$$

and the number of non-zero coordinates in this vector is greater than that in $A u$. We may now iterate this process, obtaining at each stage at least one more non-zero coordinate than we had before. We eventually obtain a vector $v>0$ such that $A v>0$. Thus $A$ has finite type.

We next identify the second case in Lemma 15.6 with that of an affine GCM.

Proposition 15.8 Let $A$ be an indecomposable GCM. Then the following conditions are equivalent:
(i) A has affine type
(ii) $\{u ; u \geq 0$ and $A u \geq 0\} \neq\{0\}, K_{A}=\{u ; A u=0\}$, and $K_{A}$ is a l-dimensional subspace of $\mathbb{R}^{n}$.

Proof. Suppose $A$ has affine type. Then there exists $u>0$ with $A u=0$. It follows that $\{u ; u \geq 0$ and $A u \geq 0\} \neq\{0\}$. Also $\lambda u \in K_{A}$ for all $\lambda \in \mathbb{R}$. By Lemma 15.6 we see that $K_{A}=\{w ; A w=0\}$ and that $K_{A}$ is a 1-dimensional subspace of $\mathbb{R}^{n}$.

Conversely suppose the three conditions of (ii) are satisfied.

$$
\text { Then corank } A=1 \text {. }
$$

Also there exists $u \neq 0$ with $u \geq 0$ and $A u \geq 0$. By Lemma 15.5 we have $u>0$. So there exists $u>0$ with $A u \geq 0$. But $K_{A}=\{u ; A u=0\}$. Hence there exists $u>0$ with $A u=0$. Finally $A u \geq 0$ implies $A u=0$. Thus $A$ has affine type.

Proposition 15.9 Let A be an indecomposable GCM. Then:

> if A has finite type $A^{\mathrm{t}}$ has finite type if A has affine type $A^{\mathrm{t}}$ has affine type.

Proof. To prove these results we shall make use of Proposition 15.3.
Suppose $A$ has finite type. We show there does not exist $v>0$ with $A v<0$. For if $A v<0$ then $A(-v)>0$ and so $-v>0$ or $-v=0$. Hence $v<0$ or $v=0$. This contradicts $v>0$. We may now apply Proposition 15.3 to show there exists $u \neq 0$ with $u \geq 0$ and $A^{\mathrm{t}} u \geq 0$. So $\left\{u ; u \geq 0\right.$ and $\left.A^{\mathrm{t}} u \geq 0\right\} \neq\{0\}$. By Lemma 15.6 either

$$
K_{A^{t}} \subset\{u ; u>0\} \cup\{0\}
$$

or $K_{A^{\mathrm{t}}}=\left\{u ; A^{\mathrm{t}} u=0\right\}$ and this is a 1-dimensional subspace. Now $\operatorname{det} A \neq 0$ so $\operatorname{det} A^{\mathrm{t}} \neq 0$. Thus the latter case cannot occur. The former case must therefore occur, so by Proposition $15.7 A^{t}$ has finite type.

Now suppose $A$ has affine type. We again show there does not exist $v>0$ with $A v<0$. For $A(-v)>0$ is impossible in the affine case.

By Proposition 15.3 there exists $u \neq 0$ with $u \geq 0$ and $A^{\mathrm{t}} u \geq 0$. So $\{u ; u \geq 0$ and $\left.A^{\mathrm{t}} u \geq 0\right\} \neq\{0\}$. By Lemma 15.6 we may again conclude that either

$$
\begin{aligned}
& K_{A^{\mathrm{t}}} \subset\{u ; u>0\} \cup\{0\} \text { or } \\
& K_{A^{\mathrm{t}}}=\{u ; A u=0\} \text { and this is a 1-dimensional subspace. }
\end{aligned}
$$

Now corank $A=1$ so corank $A^{\mathrm{t}}=1$. This shows that we cannot have the first possibility. Thus the second possibility holds, and then by Proposition 15.8 we see that $A^{\mathrm{t}}$ has affine type.

We may now identify the case not appearing in Lemma 15.6 with that of an indefinite GCM.

Proposition 15.10 Let $A$ be an indecomposable GCM. Then the following conditions are equivalent:

> A has indefinite type $\{u ; u \geq 0$ and $A u \geq 0\}=\{0\}$.

Proof. Suppose $A$ has indefinite type. Then $u \geq 0$ and $A u \geq 0$ imply $u=0$.
Conversely suppose $\{u ; u \geq 0$ and $A u \geq 0\}=\{0\}$. Then the same condition holds for $A^{\mathrm{t}}$, i.e. $\left\{u ; u \geq 0\right.$ and $\left.A^{\mathrm{t}} u \geq 0\right\}=\{0\}$. This follows from Lemma 15.6 and Propositions $15.7,15.8$ and 15.9. But then Proposition 15.3 shows that there exists $v>0$ with $A v<0$. Thus $A$ has indefinite type.

We are now able to achieve our aim of proving Theorem 15.1. For each indecomposable GCM A Lemma 15.6 shows that exactly one of the following conditions holds:
(a) $\{u ; u \geq 0$ and $A u \geq 0\} \neq\{0\}$ and $K_{A} \subset\{u ; u>0\} \cup\{0\}$.
(b) $\{u ; u \geq 0$ and $A u \geq 0\} \neq\{0\}, K_{A}=\{u ; A u=0\}$, and $K_{A}$ is a 1-dimensional subspace.
(c) $\{u ; u \geq 0$ and $A u \geq 0\}=\{0\}$.

By Proposition $15.7 A$ satisfies (a) if and only if $A$ has finite type. By Proposition $15.8 A$ satisfies (b) if and only if $A$ has affine type. By Proposition $15.10 A$ satisfies (c) if and only if $A$ has indefinite type. Thus we have the required trichotomy for GCMs. Moreover Proposition 15.9 shows
that the type of $A^{t}$ is the same as the type of $A$. This completes the proof of Theorem 15.1

Corollary 15.11 Let $A$ be an indecomposable GCM. Then:
(a) A has finite type if and only if there exists $u>0$ with $A u>0$.
(b) A has affine type if and only if there exists $u>0$ with $A u=0$.
(c) A has indefinite type if and only if there exists $u>0$ with $A u<0$.

Proof. (a) Suppose $u>0$ and $A u>0$. A cannot have affine type as then $A u \geq 0$ would imply $A u=0$. $A$ cannot have indefinite type as then $u \geq 0$ and $A u \geq 0$ would imply $u=0$. Thus $A$ has finite type.
(b) Suppose $u>0$ and $A u=0$. $A$ cannot have finite type since $\operatorname{det} A=0$. $A$ cannot have indefinite type since then $u \geq 0$ and $A u \geq 0$ would imply $u=0$. Thus $A$ has affine type.
(c) Suppose $u>0$ and $A u<0$. Then $A(-u)>0$. A cannot have finite type as this would imply $-u>0$ or $-u=0$. A cannot have affine type since $A(-u)>0$ would then imply $A(-u)=0$. Thus $A$ has indefinite type.

Remark 15.12 In proving the results of Section 15.1 we have assumed that $A$ is a GCM. However, we have not used the full force of this assumption. Inspection of the proofs shows that we have nowhere assumed that $A_{i i}=2$ or that $A_{i j} \in \mathbb{Z}$. This remark will be useful in some subsequent applications.

### 15.2 Symmetrisable generalised Cartan matrices

In this section we shall consider a special type of GCM which plays a key role in the theory of Kac-Moody algebras. These are the symmetrisable GCMs. Before giving the definition we obtain some preliminary results.

Let $A=\left(A_{i j}\right)$ be a GCM with $i, j \in\{1, \ldots, n\}$ and let $J$ be a subset of $\{1, \ldots, n\}$. Let $A_{J}=\left(A_{i j}\right), i, j \in J$. Then $A_{J}$ is also a GCM, called a principal minor of $A$.

Lemma 15.13 (i) Suppose $A$ is an indecomposable GCM of finite type and $A_{J}$ is an indecomposable principal minor of $A$. Then $A_{J}$ also has finite type.
(ii) Suppose $A$ is an indecomposable GCM of affine type and $A_{J}$ is a proper indecomposable principal minor of $A$. Then $A_{J}$ has finite type.

Proof. (i) By passing to an equivalent GCM we may choose the numbering so that $J=\{1, \ldots, m\}$ for some $m \leq n$. Let $K=\{m+1, \ldots, n\}$. Let

$$
\begin{gathered}
A=\left(\begin{array}{ll}
P & Q \\
R & S
\end{array}\right)_{n-m}^{m} \\
m n-m
\end{gathered}
$$

Now there exists $u>0$ with $A u>0$. Let $u=\binom{u_{J}}{u_{K}}$. Then

$$
A u=\left(\begin{array}{cc}
P & Q \\
R & S
\end{array}\right)\binom{u_{J}}{u_{K}}=\binom{P u_{J}+Q u_{K}}{R u_{J}+S u_{K}}
$$

Since $A u>0$ we have $P u_{J}+Q u_{K}>0$. However, $Q u_{K} \leq 0$ so $P u_{J}>0$. Thus there exists $u_{J}>0$ with $A_{J} u_{J}>0$. By Corollary $15.11 A_{J}$ has finite type.
(ii) As before we may assume $J=\{1, \ldots, m\}$. This time we have $m<n$. Let

$$
\begin{aligned}
A= & \left(\begin{array}{ll}
P & Q \\
R & S
\end{array}\right)_{n-m}^{m} \quad \text { where } P=A_{J} \\
& m n-m
\end{aligned}
$$

Since $A$ has affine type there exists $u>0$ with $A u=0$. We have

$$
A u=\left(\begin{array}{cc}
P & Q \\
R & S
\end{array}\right)\binom{u_{J}}{u_{K}}=\binom{P u_{J}+Q u_{K}}{R u_{J}+S u_{K}} .
$$

Hence $P u_{J}+Q u_{K}=0$. Now $Q u_{K} \leq 0$ so $P u_{J} \geq 0$.
Suppose if possible that $P u_{J}=0$. Then $Q u_{K}=0$, and since $u_{K}>0$ this implies that $Q=O$. But then $R=O$ also and $A$ is decomposable, a contradiction. Hence we have $u_{J}>0, P u_{J} \geq 0, P u_{J} \neq 0$. This implies that $P=A_{J}$ cannot have affine type or indefinite type. Thus $A_{J}$ has finite type.

We next describe our trichotomy in the special case in which the indecomposable GCM is symmetric.

Proposition 15.14 Suppose $A$ is a symmetric indecomposable GCM. Then:
(a) A has finite type if and only if $A$ is positive definite.
(b) A has affine type if and only if $A$ is positive semidefinite of corank 1.
(c) A has indefinite type if and only if A satisfies neither of these conditions.

Proof. (a) Let $A$ have finite type. Then there exists $u>0$ with $A u>0$. Hence for all $\lambda \geq 0$ we have $(A+\lambda I) u>0$. Thus $A+\lambda I$ has finite type by Corollary 15.11. (Note that $A+\lambda I$ need not be a GCM, but the results of Section 15.1 can be applied to it by Remark 15.12.) Thus $\operatorname{det}(A+\lambda I) \neq 0$
when $\lambda \geq 0$, that is $\operatorname{det}(A-\lambda I) \neq 0$ when $\lambda \leq 0$. Now the eigenvalues of the real symmetric matrix $A$ are all real. Thus all the eigenvalues of $A$ must be positive. Hence $A$ is positive definite.

Conversely suppose $A$ is positive definite. Then $\operatorname{det} A \neq 0$ so $A$ has finite or indefinite type. If $A$ has indefinite type there exists $u>0$ with $A u<0$. But then $u^{\mathrm{t}} A u<0$, contradicting the fact that $A$ is positive definite. Thus $A$ must have finite type.
(b) Let $A$ have affine type. Then there exists $u>0$ with $A u=0$. Hence for all $\lambda>0$ we have $(A+\lambda I) u>0$. Thus by Corollary $15.11 A+\lambda I$ has finite type when $\lambda>0$. (We are again using Remark 15.12 here.) Thus $\operatorname{det}(A+\lambda I) \neq 0$ when $\lambda>0$, that is $\operatorname{det}(A-\lambda I) \neq 0$ when $\lambda<0$. Thus all eigenvalues of $A$ are non-negative. But $A$ has corank 1 so 0 occurs as an eigenvalue with multiplicity 1 , and the remaining eigenvalues are all positive. Hence $A$ is positive semi-definite of corank 1.

Conversely suppose $A$ is positive semi-definite of corank 1. Then $\operatorname{det} A=0$ so $A$ cannot have finite type. Suppose $A$ has indefinite type. Then there exists $u>0$ with $A u<0$. Thus $u^{\mathrm{t}} A u<0$, which contradicts the fact that $A$ is positive semi-definite. Thus $A$ must have affine type.
(c) This follows from (a) and (b).

In general a GCM need not be symmetric, but it may nevertheless satisfy the weaker condition of being symmetrisable.

Definition A GCM A is symmetrisable if there exists a non-singular diagonal matrix $D$ and a symmetric matrix $B$ such that $A=D B$.

Lemma 15.15 Let A be a GCM. Then A is symmetrisable if and only if

$$
A_{i_{1} i_{2}} A_{i_{2} i_{3}} \ldots A_{i_{k} i_{1}}=A_{i_{2} i_{1}} A_{i_{3} i_{2}} \ldots A_{i_{1} i_{k}}
$$

for all $i_{1}, i_{2}, \ldots, i_{k} \in\{1, \ldots, n\}$.

Proof. Suppose $A$ is symmetrisable. Then $A=D B$ with $D=\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right)$ and $B=\left(B_{i j}\right)$. Thus $A_{i j}=d_{i} B_{i j}$. Hence

$$
\begin{aligned}
A_{i_{1} i_{2}} \ldots A_{i_{k} i_{1}} & =d_{i_{1}} \ldots d_{i_{k}} B_{i_{1} i_{2}} \ldots B_{i_{k} i_{1}} \\
A_{i_{2} i_{1}} \ldots A_{i_{1} i_{k}} & =d_{i_{1}} \ldots d_{i_{k}} B_{i_{2} i_{1}} \ldots B_{i_{1} i_{k}}
\end{aligned}
$$

and these are equal since $B$ is symmetric.
Conversely suppose

$$
A_{i_{1} i_{2}} \ldots A_{i_{k} i_{1}}=A_{i_{2} i_{1}} \ldots A_{i_{1} i_{k}}
$$

for all $i_{1}, \ldots, i_{k}$. We may suppose $A$ is indecomposable since the result in this case implies it for all $A$. Thus for each $i \in\{1, \ldots, n\}$ there exists a sequence

$$
1=j_{1}, j_{2}, \ldots, j_{t}=i
$$

with

$$
A_{j_{1} j_{2}} \neq 0, \quad A_{j_{2} j_{3}} \neq 0, \quad \ldots, \quad A_{j_{t-1} j_{t}} \neq 0
$$

We choose a number $d_{1} \neq 0$ in $\mathbb{R}$. We wish to define $d_{i}$ by

$$
d_{i}=\frac{A_{j_{i} j_{t-1}} \ldots A_{j_{2} j_{1}}}{A_{j_{1} j_{2}} \ldots A_{j_{t-1} j_{t}}} d_{1}
$$

However, we must check that this definition of $d_{i}$ depends only upon $i$ and not on the sequence chosen from 1 to $i$. So let

$$
1=k_{1}, k_{2}, \ldots, k_{u}=i
$$

be a second such sequence from 1 to $i$. We claim that

$$
\frac{A_{j_{t} j_{t-1}} \ldots A_{j_{2} j_{1}}}{A_{j_{1} j_{2}} \ldots A_{j_{t-1} j_{t}}}=\frac{A_{k_{u} k_{u-1}} \ldots A_{k_{2} k_{1}}}{A_{k_{1} k_{2}} \ldots A_{k_{u-1} k_{u}}}
$$

that is $A_{1 k_{2}} A_{k_{2} k_{3}} \ldots A_{k_{u-1} i} A_{i j_{t-1}} \ldots A_{j_{2} 1}=A_{k_{2} 1} A_{k_{3} k_{2}} \ldots A_{i k_{u-1}} A_{j_{t-1} i} \ldots A_{1 j_{2}}$. This is in fact one of the given conditions on the matrix $A$. Thus $d_{i} \in \mathbb{R}$ is well defined and $d_{i} \neq 0$. Let $D=\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right)$. Define $B_{i j}$ by $A_{i j}=d_{i} B_{i j}$. We show that $B_{j i}=B_{i j}$, that is $\frac{A_{j i}}{d_{j}}=\frac{A_{i j}}{d_{i}}$. If $A_{i j}=0$ then $A_{j i}=0$ also and the condition is satisfied. So suppose $A_{i j} \neq 0$. Let $1=j_{1}, j_{2}, \ldots, j_{t}=i$ be a sequence from 1 to $i$ of the type described above. Then $1=j_{1}, j_{2}, \ldots, j_{t}, j$ is such a sequence from 1 to $j$. These sequences may be used to obtain $d_{i}$ and $d_{j}$ respectively, and we have

$$
d_{j}=\frac{A_{j i}}{A_{i j}} d_{i}
$$

Thus $B_{j i}=B_{i j}$. Hence $A=D B$ where $D$ is diagonal and non-singular, and $B$ is symmetric. Thus $A$ is symmetrisable.

Corollary 15.16 Let $A$ be a symmetrisable indecomposable GCM. Then $A$ can be expressed in the form $A=D B$ where $D=\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right), B$ is symmetric, with $d_{1}, \ldots, d_{n}>0$ in $\mathbb{Z}$ and $B_{i j} \in \mathbb{Q}$. Also $D$ is determined by these conditions up to a scalar multiple.

Proof. We choose any $d_{1} \in \mathbb{Q}$ with $d_{1}>0$. Then Lemma 15.15 shows that $d_{i} \in \mathbb{Q}$ and $d_{i}>0$ for each $i$. Thus by multiplying by a positive scalar
we may assume each $d_{i} \in \mathbb{Z}$ with $d_{i}>0$. Also $B_{i j}=A_{i j} / d_{i}$ lies in $\mathbb{Q}$. The proof of Lemma 15.15 also shows that $D$ is determined up to a scalar multiple.

The following important result shows that indecomposable GCMs in the first two classes of our trichotomy are symmetrisable.

Theorem 15.17 Let A be an indecomposable GCM of finite or affine type. Then $A$ is symmetrisable.

Proof. First suppose there is no set of integers $i_{1}, i_{2}, \ldots, i_{k}$ with $k \geq 3$ such that $i_{1} \neq i_{2}, i_{2} \neq i_{3}, \ldots, i_{k-1} \neq i_{k}, i_{k} \neq i_{1}$ and

$$
A_{i_{1} i_{2}} \neq 0, A_{i_{2} i_{3}} \neq 0, \ldots, A_{i_{k-1} i_{k}} \neq 0, A_{i_{k} i_{1}} \neq 0
$$

Then Lemma 15.15 shows that $A$ is symmetrisable.
Thus we suppose there is such a sequence $i_{1}, \ldots, i_{k}$ with $k \geq 3$ and we choose such a sequence with minimal possible value of $k$. We thus have

$$
A_{i_{r} i_{s}} \neq 0 \quad \text { if }(r, s) \in\{(1,2),(2,3), \ldots,(k, 1),(2,1),(3,2), \ldots,(1, k)\}
$$

The minimality of $k$ shows that $A_{i_{r} i_{s}}=0$ if $(r, s)$ does not lie in the above set. Otherwise there would be such a sequence with a smaller value of $k$.

Let $J=\left\{i_{1}, \ldots, i_{k}\right\}$. Then the principal minor $A_{J}$ of $A$ has form

$$
A_{J}=\left(\begin{array}{cccccccc}
2 & -r_{1} & 0 & \cdot & \cdot & . & 0 & -s_{k} \\
-s_{1} & 2 & -r_{2} & \cdot & & & 0 \\
0 & -s_{2} & 2 & \cdot & . & & \cdot \\
\cdot & \cdot & \cdot & \cdot & \cdot & \cdot & & \cdot \\
\cdot & & \cdot & \cdot & \cdot & \cdot & \cdot & \cdot \\
\cdot & & & \cdot & 2 & \cdot & 0 \\
0 & & & \cdot & \cdot & \cdot & 2 & -r_{k-1} \\
-r_{k} & 0 & . & \cdot & \cdot & 0 & -s_{k-1} & 2
\end{array}\right) .
$$

with $r_{i}, s_{i} \in \mathbb{Z}$ satisfying $r_{i}>0, s_{i}>0$. In particular we see that $A_{J}$ is indecomposable. Now $A_{J}$ must have finite or affine type by Lemma 15.13. Thus there exists $u>0$ with $A_{J} u \geq 0$. Let $u=\left(u_{1}, \ldots, u_{k}\right)$. We define the $k \times k$ matrix $M$ by

$$
M=\operatorname{diag}\left(u_{1}^{-1}, \ldots, u_{k}^{-1}\right) A_{J} \operatorname{diag}\left(u_{1}, \ldots, u_{k}\right)
$$

Then $M_{i j}=u_{i}^{-1}\left(A_{J}\right)_{i j} u_{j}$. Thus

$$
\sum_{j} M_{i j}=u_{i}^{-1} \sum_{j}\left(A_{J}\right)_{i j} u_{j} \geq 0
$$

In particular we have $\sum_{i, j} M_{i j} \geq 0$. Now we have

$$
M=\left(\begin{array}{cccccccc}
2 & -r_{1}^{\prime} & 0 & \cdot & \cdot & \cdot & 0 & -s_{k}^{\prime} \\
-s_{1}^{\prime} & 2 & -r_{2}^{\prime} & \cdot & & & & 0 \\
0 & -s_{2}^{\prime} & 2 & \cdot & \cdot & & \cdot \\
\cdot & \cdot & \cdot & \cdot & \cdot & \cdot & & \cdot \\
\cdot & \cdot & \cdot & \cdot & \cdot & \cdot & \cdot & \cdot \\
\cdot & & \cdot & \cdot & \cdot & \cdot & 0 \\
0 & & & & \cdot & \cdot & 2 & -r_{k-1}^{\prime} \\
-r_{k}^{\prime} & 0 & . & \cdot & \cdot & 0 & -s_{k-1}^{\prime} & 2
\end{array}\right) .
$$

with $r_{i}^{\prime}=u_{i}^{-1} r_{i} u_{i+1}, s_{i}^{\prime}=u_{i+1}^{-1} s_{i} u_{i}$. (We define $u_{k+1}=u_{1}$. )
We note that $r_{i}^{\prime}>0, s_{i}^{\prime}>0$ and $r_{i}^{\prime} s_{i}^{\prime}=r_{i} s_{i} \in \mathbb{Z}$. We also have

$$
\sum_{i, j} M_{i j}=2 k-\left(r_{1}^{\prime}+s_{1}^{\prime}\right)-\cdots-\left(r_{k}^{\prime}+s_{k}^{\prime}\right) .
$$

Now $\frac{r_{i}^{\prime}+s_{i}^{\prime}}{2} \geq \sqrt{r_{i}^{\prime} s_{i}^{\prime}}=\sqrt{r_{i} s_{i}} \geq 1$ hence $r_{i}^{\prime}+s_{i}^{\prime} \geq 2$. Since $\sum_{i, j} M_{i j} \geq 0$ we deduce that $r_{i}^{\prime}+s_{i}^{\prime}=2$ and $r_{i}^{\prime} s_{i}^{\prime}=1$. Hence $r_{i} s_{i}=1$ and, since $r_{i}$, $s_{i}$ are positive integers, we have $r_{i}=1, s_{i}=1$. It follows that

$$
A_{J}=\left(\begin{array}{cccccccc}
2 & -1 & 0 & . & . & . & 0 & -1 \\
-1 & 2 & -1 & . & & & & 0 \\
0 & -1 & 2 & . & . & & & . \\
. & . & . & . & . & . & & . \\
. & & . & . & . & . & . & . \\
. & & & . & . & . & . & 0 \\
0 & & & & . & . & 2 & -1 \\
-1 & 0 & . & . & . & 0 & -1 & 2
\end{array}\right)
$$

Let $v=(1, \ldots, 1)$. Then $v>0$ and $A_{J} v=0$. Thus $A_{J}$ has affine type by Corollary 15.11 . Lemma 15.13 shows that this can only happen when $A_{J}=A$. Thus $A$ is symmetric, in particular symmetrisable as required.

We are now able to prove the following basic description of our trichotomy. It generalises the description previously obtained in Proposition 15.14 for symmetric indecomposable GCMs.

Theorem 15.18 Let $A$ be an indecomposable GCM. Then:
(a) A has finite type if and only if all its principal minors have positive determinant.
(b) A has affine type if and only if $\operatorname{det} A=0$ and all proper principal minors have positive determinant.
(c) A has indefinite type if and only if A satisfies neither of these two conditions.

Proof. (a) Suppose $A$ has finite type. Then $A$ is symmetrisable by Theorem 15.17, hence $A=D B$ where $D=\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right)$ with $d_{i}>0$ and $B$ is symmetric, by Corollary 15.16. The matrix $B$ need not necessarily be a GCM, but Remark 15.12 shows that we can nevertheless define the type of $B$. Moreover Corollary 15.11 shows that $A$ and $B$ have the same type. Thus $B$ is a symmetric indecomposable matrix of finite type, and so $\operatorname{det} B>0$ by Proposition 15.14. It follows that $\operatorname{det} A>0$ also. Now all principal minors of $A$ also have finite type by Lemma 15.13. Thus these also have positive determinant.

Conversely suppose that all principal minors of $A$ have positive determinant. Suppose there is a set of integers $i_{1}, \ldots, i_{k}$ with $k \geq 3$ such that $i_{1} \neq i_{2}, i_{2} \neq i_{3}, \ldots, i_{k} \neq i_{1}$ with

$$
A_{i_{1} i_{2}} A_{i_{2} i_{3}} \ldots A_{i_{k} i_{1}} \neq 0
$$

Choose such a sequence with minimal possible $k$ and let $J=\left\{i_{1}, \ldots, i_{k}\right\}$. Then the principal minor $A_{J}$ of $A$ has form

$$
A_{J}=\left(\begin{array}{cccccccc}
2 & -1 & 0 & \cdot & \cdot & \cdot & 0 & -1 \\
-1 & 2 & -1 & \cdot & & & 0 \\
0 & -1 & 2 & \cdot & \cdot & & \cdot \\
\cdot & \cdot & \cdot & \cdot & \cdot & \cdot & & \cdot \\
\cdot & & \cdot & \cdot & \cdot & \cdot & \cdot & \cdot \\
\cdot & & & \cdot & \cdot & \cdot & \cdot & 0 \\
0 & & & & . & 2 & -1 \\
-1 & 0 & \cdot & \cdot & \cdot & 0 & -1 & 2
\end{array}\right)
$$

by the proof of Theorem 15.17. But then $\operatorname{det} A_{J}=0$, a contradiction. Thus there is no such sequence $i_{1}, \ldots, i_{k}$. By Lemma $15.15 A$ is symmetrisable. Hence $A=D B$ where $D=\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right)$ with $d_{i}>0$ and $B$ is symmetric. Again $B$ need not be a GCM but we can nevertheless define its type using Remark 15.12 and, by Corollary $15.11, A$ and $B$ have the same type. Now the principal minors of the symmetric matrix $B$ all have positive determinant
and so $B$ is positive definite. Thus $B$ has finite type by Proposition 15.14 , and so $A$ has finite type also.
(b) Now suppose $A$ has affine type. Then $\operatorname{det} A=0$. All proper principal minors of $A$ have finite type by Lemma 15.13 and so have positive determinant by (a).

Suppose conversely that $\operatorname{det} A=0$ and that all proper principal minors of $A$ have positive determinant. Suppose there is a set of integers $i_{1}, \ldots, i_{k}$ with $k \geq 3$ such that $i_{1} \neq i_{2}, \ldots, i_{k} \neq i_{1}$ with

$$
A_{i_{1} i_{2}} A_{i_{2} i_{3}} \ldots A_{i_{k} i_{1}} \neq 0
$$

Choose such a sequence with minimal $k$, and let $J=\left\{i_{1}, \ldots, i_{k}\right\}$. Then the principal minor $A_{J}$ has form

$$
A_{J}=\left(\begin{array}{cccccccc}
2 & -1 & 0 & \cdot & \cdot & \cdot & 0 & -1 \\
-1 & 2 & -1 & \cdot & & & & 0 \\
0 & -1 & 2 & \cdot & \cdot & & & \cdot \\
\cdot & \cdot & \cdot & \cdot & \cdot & \cdot & & \cdot \\
\cdot & & \cdot & \cdot & \cdot & \cdot & \cdot & \cdot \\
\cdot & & & \cdot & \cdot & \cdot & \cdot & 0 \\
0 & & & & \cdot & \cdot & 2 & -1 \\
-1 & 0 & \cdot & \cdot & \cdot & 0 & -1 & 2
\end{array}\right) .
$$

as above. Since $\operatorname{det} A_{J}=0$ we have $A_{J}=A$. But then $A$ is affine since $A u=0$ with $u=(1, \ldots, 1)$. Thus suppose there is no such sequence $i_{1}, \ldots, i_{k}$. Then $A$ is symmetrisable by Lemma 15.15 , and has form $A=D B$ where $D=\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right)$ with $d_{i}>0$ and $B$ is symmetric. Now $\operatorname{det} B=0$ and all proper principal minors of $B$ have positive determinant. This implies that the symmetric matrix $B$ is positive semidefinite of corank 1. Hence $B$ is of affine type by Proposition 15.14 . Thus $A$ has affine type also, by Corollary 15.11.
(c) This follows directly from (a) and (b).

### 15.3 The classification of affine generalised Cartan matrices

In this section we shall determine explicitly which indecomposable GCMs lie in each class of our trichotomy. We begin with indecomposable GCMs of finite type.

Theorem 15.19 Let $A$ be an indecomposable GCM. Then A has finite type if and only if $A$ is a Cartan matrix. Thus the indecomposable GCMs of finite type are those on the standard list 6.12.

Proof. We recall from Sections 6.1 and 6.2 that a GCM is a Cartan matrix if and only if it satisfies the conditions:
(a) $A_{i j} \in\{0,-1,-2,-3\}$ for all $i \neq j$
(b) $A_{i j}=-2$ or -3 implies $A_{j i}=-1$
(c) the quadratic form

$$
\mathrm{Q}\left(x_{1}, \ldots, x_{n}\right)=2 \sum_{i=1}^{n} x_{i}^{2}-\sum_{i \neq j} \sqrt{ } n_{i j} x_{i} x_{j}
$$

is positive definite, where $n_{i j}=A_{i j} A_{j i}$.
Suppose $A$ is a Cartan matrix. Then $A_{i j}=2 \frac{\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}$. Let $D=\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right)$ where $d_{i}=\sqrt{ }\left\langle\alpha_{i}, \alpha_{i}\right\rangle$. Then $\left(D A D^{-1}\right)_{i j}=2 \frac{\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\sqrt{ }\left\langle\alpha_{i}, \alpha_{i}\right\rangle \sqrt{ }\left\langle\alpha_{j}, \alpha_{j}\right\rangle}$ and so $D A D^{-1}$ is the matrix of the quadratic form $\mathrm{Q}\left(x_{1}, \ldots, x_{n}\right)$. Since Q is positive definite $\operatorname{det}\left(D A D^{-1}\right)>0$ and so $\operatorname{det} A>0$.

Now any principal minor $A_{J}$ of the Cartan matrix $A$ is also a Cartan matrix. Hence $\operatorname{det} A_{J}>0$ for all principal minors of $A$. Thus $A$ has finite type by Theorem 15.18 (a).

Now suppose conversely that $A$ has finite type. Suppose $i \neq j$ and consider the $2 \times 2$ principal minor

$$
\left(\begin{array}{cc}
2 & A_{i j} \\
A_{j i} & 2
\end{array}\right)
$$

By Theorem 15.18 (a) the determinant of this minor is positive, hence $A_{i j} A_{j i}<4$. Since $A_{i j}$ and $A_{j i}$ are both non-positive integers such that $A_{i j}=0$ if and only if $A_{j i}=0$ we deduce that $A_{i j} \in\{0,-1,-2,-3\}$ and that $A_{i j} \in$ $\{-2,-3\}$ implies $A_{j i}=-1$.

Since $A$ has finite type $A$ is symmetrisable by Theorem 15.17. Thus $A=$ $D B$ where $D=\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right), d_{i}>0$, and $B$ is symmetric. Although $B$ need not be a GCM we may define the type of $B$ by using Remark 15.12. Thus $B$ is indecomposable of finite type, and so $B$ is positive definite by Proposition 15.14 (a). Let $y_{i}=\sqrt{ } d_{i} x_{i}$. Then

$$
\mathrm{Q}\left(x_{1}, \ldots, x_{n}\right)=2 \sum_{i} x_{i}^{2}-\sum_{i \neq j} \sqrt{ } n_{i j} x_{i} x_{j}
$$

$$
\begin{aligned}
& =\frac{2}{d_{i}} \sum_{i} y_{i}^{2}-\sum_{i \neq j} \frac{1}{\sqrt{ } d_{i}} \frac{1}{\sqrt{ } d_{j}} \sqrt{ }\left(A_{i j} A_{j i}\right) y_{i} y_{j} \\
& =\sum_{i} B_{i i} y_{i}^{2}+\sum_{i \neq j} B_{i j} y_{i} y_{j}
\end{aligned}
$$

Since $B$ is positive definite we see that $\mathrm{Q}\left(x_{1}, \ldots, x_{n}\right)$ is positive definite. Thus $A$ is a Cartan matrix.

Having determined the indecomposable GCMs of finite type we next determine those of affine type. This will also determine those of indefinite type, as those remaining.

To each GCM $A$ we define an associated diagram $\Delta(A)$ called the Dynkin diagram of $A$. This extends the definition of the Dynkin diagram of a Cartan matrix given in Section 6.2. The vertices of $\Delta(A)$ are labelled $1, \ldots, n$ where $A$ is an $n \times n$ matrix. Suppose $i, j$ are distinct vertices of $\Delta(A)$. We explain how $i, j$ are joined in $\Delta(A)$. This depends on the pair $\left(A_{i j}, A_{j i}\right)$. We recall that $A_{i j}$ and $A_{j i}$ lie in $\mathbb{Z}, A_{i j} \leq 0, A_{j i} \leq 0$ and $A_{i j}=0$ if and only if $A_{j i}=0$. The rules are as follows.
(a) If $A_{i j} A_{j i}=0$ vertices $i, j$ are not joined.
(b) If $A_{i j} A_{j i}=1$ vertices $i, j$ are joined by a single edge.
(c) If $A_{i j} A_{j i}=2, \quad A_{i j}=-1, \quad A_{j i}=-2$ vertices $i, j$ are joined by a double edge with an arrow pointing towards $j$.
(d) If $A_{i j} A_{j i}=3, \quad A_{i j}=-1, \quad A_{j i}=-3$ vertices $i, j$ are joined by a triple edge with an arrow pointing towards $j$.
(e) If $A_{i j} A_{j i}=4, \quad A_{i j}=-1, \quad A_{j i}=-4$ vertices $i, j$ are joined by a quadruple edge with an arrow pointing towards $j$.
(f) If $A_{i j} A_{j i}=4, \quad A_{i j}=-2, \quad A_{j i}=-2$ vertices $i, j$ are joined by a double edge with two arrows pointing away from $i, j$.

(g) If $A_{i j} A_{j i} \geq 5$ vertices $i, j$ are joined by an edge with the numbers $\left|A_{i j}\right|,\left|A_{j i}\right|$ shown on it.


It is clear that the GCM $A$ is determined by its Dynkin diagram $\Delta(A)$. Moreover $A$ is indecomposable if and only if $\Delta(A)$ is connected.

We now consider a set of connected Dynkin diagrams called the affine list.
15.20 The affine list of Dynkin diagrams
$\sigma><0$
$\tilde{A}_{1}$



$\tilde{B}_{3}^{1}$




$\tilde{C}_{3}$

$\tilde{C}_{4}$


$\tilde{C}_{3}^{\mathrm{t}}$
$\sigma<0-0-\sigma>0 \cdots$
$\tilde{C}_{4}^{\mathrm{t}}$


$\tilde{D}_{5}$



$\tilde{E}_{7}$



We note that many, but not all, of the Dynkin diagrams on the affine list appeared in Lemma 6.8. We also note that every proper connected subdiagram of a Dynkin diagram on the affine list appears on list 6.11 of Dynkin diagrams of finite type. We shall call this the finite list.

Proposition 15.21 Let A be a GCM whose Dynkin diagram lies on the affine list. Then $\operatorname{det} A=0$.

Proof. First suppose that $\Delta(A)$ has 2 vertices. Then either $\Delta(A)=\tilde{A}_{1}$ and $A=\left(\begin{array}{cc}2 & -2 \\ -2 & 2\end{array}\right)$ or $\Delta(A)=\tilde{A}_{1}^{\prime}$ and $A=\left(\begin{array}{cc}2 & -4 \\ -1 & 2\end{array}\right)$. In either case $\operatorname{det} A=0$.

Next suppose that $\Delta(A)=\tilde{A}_{l}$ for $l \geq 2$. Then the sum of all the rows of $A$ is zero, and so $\operatorname{det} A=0$.

In all other cases $\Delta(A)$ has a vertex, say 1 , joined to just one other vertex, say 2 . Moreover we can choose these vertices so that they are joined by a single or a double edge. In the case of a single edge we have

$$
\operatorname{det} A=2 \operatorname{det} B-\operatorname{det} C
$$

where $B$ is obtained from $A$ by removing row and column 1 , and $C$ is obtained from $B$ by removing row and column 2 . This relation between determinants is obtained as in the proof of Theorem 6.7. The connected components of $B$ and $C$ are Cartan matrices of finite type, so their determinants are known from the proof of Theorem 6.7. In all cases this gives $\operatorname{det} A=0$.

In the case when vertices 1,2 are joined by a double edge we obtain

$$
\operatorname{det} A=2 \operatorname{det} B-2 \operatorname{det} C
$$

again as in the proof of Theorem 6.7. Again $B, C$ have connected components of finite type so we know their determinants, and in each case we obtain $\operatorname{det} A=0$.

Proposition 15.22 Let A be a GCM whose Dynkin diagram lies on the affine list. Then A has affine type.

Proof. By Proposition 15.21 we have $\operatorname{det} A=0$. Also the Dynkin diagram of any proper principal minor has connected components on the finite list. Thus all proper principal minors have positive determinant. It follows that $A$ has affine type by Theorem 15.18 (b).

We shall now prove the converse.

Theorem 15.23 Let A be an indecomposable GCM. Then A has affine type if and only if its Dynkin diagram $\Delta(A)$ lies on the affine list.

Proof. Suppose $A$ has affine type. Then every proper indecomposable principal minor of $A$ has finite type, by Lemma 15.13 (ii). Thus all proper connected subdiagrams of $\Delta(A)$ lie on the finite list, by Theorem 15.19.

If $\Delta(A)$ has only one vertex $A$ has finite type, so there is no possible affine $A$.

If $\Delta(A)$ has two vertices then

$$
A=\left(\begin{array}{cc}
2 & -a \\
-b & 2
\end{array}\right)
$$

where $a, b$ are positive integers. Since $\operatorname{det} A=0$ we have $a b=4$. The possibilities are $(a, b)=(1,4)(4,1)(2,2)$. Thus $\Delta(A)=\tilde{A}_{1}$ or $\tilde{A}_{1}^{\prime}$.

Now suppose $\Delta(A)$ has at least three vertices. If $\Delta(A)$ contains a cycle then the proof of Theorem 15.17 shows that $\Delta(A)=\tilde{A}_{l}$ for some $l \geq 2$. Thus we suppose that $\Delta(A)$ contains no cycle. Since all the connected subdiagrams with two vertices lie on the finite list all edges of $\Delta(A)$ have one of the forms


Suppose $\Delta(A)$ has a triple edge $\Longrightarrow$ Then $\Delta(A)$ must have exactly three vertices, otherwise $\Delta(A)$ would have a proper connected subdiagram with three vertices containing a triple edge, whereas there is no such diagram on the finite list. Thus we have

$$
A=\left(\begin{array}{ccc}
2 & -1 & 0 \\
-3 & 2 & -a \\
0 & -b & 2
\end{array}\right) \quad \text { or } \quad\left(\begin{array}{ccc}
2 & -3 & 0 \\
-1 & 2 & -a \\
0 & -b & 2
\end{array}\right)
$$

where $a, b$ are positive integers. Thus $\operatorname{det} A=2(1-a b)$. However, $\operatorname{det} A=0$ and so $a=1, b=1$. Thus $\Delta(A)=\tilde{G}_{2}$ or $\tilde{G}_{2}^{\mathrm{t}}$.

So we now suppose that $\Delta(A)$ has no triple edge. Now $\Delta(A)$ has at most two double edges, as every proper connected subdiagram appears on the finite list so has at most one. Suppose $\Delta(A)$ has two double edges.

Then every edge which can be removed to give a connected subdiagram must be a double edge. This implies that $\Delta(A)$ must be one of $\tilde{C}_{l}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\prime}$.

Thus we suppose that $\Delta(A)$ has just one double edge. If $\Delta(A)$ has a branch point then no proper connected subdiagram can contain both a double edge and a branch point, since the subdiagram lies on the finite list. This implies that $\Delta(A)$ is $\tilde{B}_{l}$ or $\tilde{B}_{l}^{\mathrm{t}}$.

Now suppose that $\Delta(A)$ has one double edge but no branch point. Then $\Delta(A)$ has form

with $a+b+2$ vertices. We have $a>0$ and $b>0$ since $\Delta(A)$ is not on the finite list. Also $b \leq 2$, otherwise there would be a proper subdiagram

and $a \leq 2$, otherwise there would be a proper subdiagram


Thus the possibilities are

$$
(a, b)=(1,1),(2,1),(1,2),(2,2)
$$

The case $(a, b)=(1,1)$ appears on the finite list so is not affine. The case $(a, b)=(2,2)$ is impossible, since it would give proper subdiagrams as above. Thus $(a, b)=(2,1)$ or $(1,2)$ and $\Delta(A)$ is $\tilde{F}_{4}$ or $\tilde{F}_{4}^{t}$.

Thus we may now assume that $\Delta(A)$ has only single edges. Consider the branch points of $\Delta(A)$. Each branch point has at most four branches, otherwise there would be a proper subdiagram

which does not appear on the finite list. If there is a branch point with four branches then $\Delta(A)=\tilde{D}_{4}$, as otherwise there would again be a proper subdiagram $\tilde{D}_{4}$.

Thus we may assume that all branch points in $\Delta(A)$ have three branches. There cannot be more than two branch points, as otherwise there would be a proper connected subdiagram with two branch points which could not be on the finite list.

Suppose $\Delta(A)$ has 2 branch points. Then any proper connected subdiagram has only one branch point, and this implies that $\Delta(A)=\tilde{D}_{l}$ for some $l \geq 5$.

So suppose $\Delta(A)$ has just one branch point. Let the branch lengths be $l_{1}, l_{2}, l_{3}$ with $l_{1} \leq l_{2} \leq l_{3}$ so that there are $l_{1}+l_{2}+l_{3}+1$ vertices. We must have $l_{1} \leq 2$, otherwise there would be a proper subdiagram

which is not on the finite list. Suppose $l_{1}=2$. Then we must have $l_{2}=2$ and $l_{3}=2$, otherwise there would again be a proper subdiagram as above. Thus $\Delta(A)=\tilde{E}_{6}$.

Thus we may assume $l_{1}=1$. Since $l_{2}=1$ would give a diagram of finite type we must have $l_{2} \geq 2$. However, $l_{2} \leq 3$ as otherwise there would be a proper subdiagram

which is not on the finite list. Thus $l_{2}=2$ or 3 .
Suppose $l_{1}=1, l_{2}=3$. Then we must have $l_{3}=3$, otherwise there would be a proper subdiagram

which is again not on the finite list. Thus $\Delta(A)=\tilde{E}_{7}$.
We may now suppose that $l_{1}=1, l_{2}=2$. Since the diagrams with $l_{3}=2,3,4$ are of finite type we must have $l_{3} \geq 5$. But $l_{3} \leq 5$ also, as otherwise there would be a proper subdiagram

which is not on the finite list. Hence $l_{3}=5$ and $\Delta(A)=\tilde{E}_{8}$.
Finally if $\Delta(A)$ has only single edges and no branch points then it lies on the finite list so $A$ cannot be affine.

Thus we have shown that whenever A is affine $\Delta(A)$ must appear on the affine list. This, together with proposition 15.22 , completes the proof.

A GCM $A$ such that $\Delta(A)$ is on the affine list will be called an affine Cartan matrix.

Corollary 15.24 Let A be an indecomposable GCM. Then A has indefinite type if and only if its Dynkin diagram $\Delta(A)$ does not appear on the finite list or the affine list.

Proof. This follows from Theorems 15.1, 15.19 and 15.23

## 16

## The invariant form, Weyl group, and root system

We now turn to the study of the Kac-Moody algebra $L(A)$ associated with a GCM $A$.

### 16.1 The invariant bilinear form

We recall from Section 4.2 that when $A$ is a Cartan matrix the corresponding finite dimensional Lie algebra $L(A)$ has a non-degenerate symmetric bilinear form

$$
\langle,\rangle: L(A) \times L(A) \rightarrow \mathbb{C}
$$

which is invariant in the sense that

$$
\langle[x y], z\rangle=\langle x,[y z]\rangle
$$

for $x, y, z \in L(A)$. The Killing form has these properties.
In the case of a GCM $A$ we cannot define the Killing form on $L(A)$ as in the finite dimensional case. We can nevertheless ask whether there is a nondegenerate, symmetric, invariant bilinear form on $L(A)$. This is not always the case, but we shall show that such a form does exist when $A$ is symmetrisable.

Thus suppose $A$ is a symmetrisable GCM. Then $A=D B$ where $D$ is diagonal and $B$ is symmetric. Let $D=\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right)$. Let $\left(H, \Pi, \Pi^{\mathrm{v}}\right)$ be a minimal realisation of $A$, where $\Pi^{\mathrm{v}}=\left\{h_{1}, \ldots, h_{n}\right\}$ is a linearly independent subset of $H, \Pi=\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ is a linearly independent subset of $H^{*}, \alpha_{j}\left(h_{i}\right)=A_{i j}$ and $\operatorname{dim} H=2 n-l$ where $l=\operatorname{rank} A$.

Let $H^{\prime}$ be the subspace of $H$ spanned by $h_{1}, \ldots, h_{n}$ and let $H^{\prime \prime}$ be a complementary subspace of $H^{\prime}$ in $H$. Then we have

$$
H=H^{\prime} \oplus H^{\prime \prime} \quad \operatorname{dim} H^{\prime}=n, \quad \operatorname{dim} H^{\prime \prime}=n-l .
$$

We define a bilinear form $\langle\rangle:, H \times H \rightarrow \mathbb{C}$ by the rules:

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=d_{i} d_{j} B_{i j} \quad i, j=1, \ldots, n \\
& \left\langle h_{i}, x\right\rangle=\left\langle x, h_{i}\right\rangle=d_{i} \alpha_{i}(x) \quad \text { for } x \in H^{\prime \prime} \\
& \langle x, y\rangle=0 \quad \text { for } x, y \in H^{\prime \prime}
\end{aligned}
$$

This is evidently a symmetric bilinear form on $H$.

## Proposition 16.1 This form on $H$ is non-degenerate.

Proof. We have $A=D B$ where $D$ is diagonal and non-singular and $B$ is symmetric. We have rank $B=l$. We observe that the symmetric matrix $B$ of rank $l$ has a non-singular $l \times l$ principal minor. If $l=n$ we can take $B$ itself as the principal minor, so suppose $l<n$. Then, for some $i$, the $i$ th row of $B$ is a linear combination of the remaining rows of $B$. Since $B$ is symmetric the $i$ th column of $B$ is a linear combination of the remaining columns of $B$. Let $B^{\prime}$ be the $(n-1) \times n$ matrix obtained from $B$ by removing the $i$ th row. Then $\operatorname{rank} B^{\prime}=l$. Let $B^{\prime \prime}$ be the $(n-1) \times(n-1)$ matrix obtained from $B^{\prime}$ by removing the $i$ th column. Then rank $B^{\prime \prime}=l$. Now $B^{\prime \prime}$ is symmetric of degree $n-1$ and rank $l$. Thus by induction $B^{\prime \prime}$ has a non-singular $l \times l$ principal minor, and this is the required principal minor of $B$.

It follows that the symmetrisable matrix $A$ has a non-singular $l \times l$ principal minor. For let $B_{J}$ be non-singular where $J$ is a subset of $\{1, \ldots, n\}$ with $|J|=l$. Then $A_{J}=D_{J} B_{J}$ where $D_{J}=\operatorname{diag}\left\{d_{j}, j \in J\right\}$ with each $d_{j} \neq 0$. Since $D_{J}$ is non-singular $A_{J}$ is also non-singular.

We now consider the special case in which $J=\{1, \ldots, l\}$. Then $A$ has form

$$
\begin{gathered}
A=\left(\begin{array}{ll}
A_{11} & A_{12} \\
A_{21} & A_{22}
\end{array}\right) \quad \begin{array}{l}
l \\
\\
\quad l \\
n-l
\end{array} \quad A_{11} \text { non-singular. }
\end{gathered}
$$

By Proposition 14.2 we may extend the linearly independent sets $h_{1}, \ldots, h_{n} \in$ $H, \alpha_{1}, \ldots, \alpha_{n} \in H^{*}$, to bases $h_{1}, \ldots, h_{2 n-l} ; \alpha_{1}, \ldots, \alpha_{2 n-l}$ such that $\alpha_{j}\left(h_{i}\right)=$ $C_{i j}$ where

$$
C=\left(\begin{array}{ccc}
A_{11} & A_{12} & O \\
A_{21} & A_{22} & I \\
O & I & O
\end{array}\right) \begin{gathered}
l \\
n-l \\
l
\end{gathered}
$$

Let

$$
D=\left(\begin{array}{cc}
D_{1} & O \\
O & D_{2}
\end{array}\right)_{n-l}^{l} .
$$

Then

$$
C=\left(\begin{array}{ccc}
D_{1} B_{11} & D_{1} B_{12} & O \\
D_{2} B_{21} & D_{2} B_{22} & I \\
O & I & O
\end{array}\right)
$$

The symmetric matrix $M$ of the bilinear form $\left\langle h_{i}, h_{j}\right\rangle \quad i, j \in\{1, \ldots, 2 n-l\}$ is

$$
M=\left(\begin{array}{ccc}
D_{1} B_{11} D_{1} & D_{1} B_{12} D_{2} & O \\
D_{2} B_{21} D_{1} & D_{2} B_{22} D_{2} & D_{2} \\
O & D_{2} & O
\end{array}\right)
$$

This matrix is non-singular since

$$
\operatorname{det} M= \pm\left(\operatorname{det} D_{1}\right)^{2}\left(\operatorname{det} D_{2}\right)^{2} \operatorname{det} B_{11} \neq 0
$$

Now suppose $A$ is any $n \times n$ symmetrisable GCM of rank $l$. Then $A$ has a non-singular $l \times l$ principal minor $A_{J}$ for some $J \subset\{1, \ldots, n\}$. Let $K$ be the complementary subset of $J$ in $\{1, \ldots, n\}$ and $L=\{n+1, \ldots, 2 n-l\}$. Then there exists a realisation $h_{1}, \ldots, h_{2 n-l} ; \alpha_{1}, \ldots, \alpha_{2 n-l}$ whose matrix $\alpha_{j}\left(h_{i}\right)=C_{i j}$ may be written symbolically in the form

$$
C=\left(\begin{array}{ccc}
A_{J} & A_{J K} & O \\
A_{K J} & A_{K} & I \\
O & I & O
\end{array}\right) \begin{gathered}
J \\
J \\
J
\end{gathered}
$$

Let

$$
D=\left(\begin{array}{cc}
D_{J} & O \\
O & D_{K}
\end{array}\right) \begin{gathered}
J \\
J
\end{gathered} \quad K
$$

Then

$$
C=\left(\begin{array}{ccc}
D_{J} B_{J} & D_{J} B_{J K} & O \\
D_{K} B_{K J} & D_{K} B_{K} & I \\
O & I & O
\end{array}\right)
$$

This time the symmetric matrix $M$ of the bilinear form $\left\langle h_{i}, h_{j}\right\rangle$ for $i, j \in$ $\{1, \ldots, 2 n-l\}$ is

$$
M=\left(\begin{array}{ccc}
D_{J} B_{J} D_{J} & D_{J} B_{J K} D_{K} & O \\
D_{K} B_{K J} D_{J} & D_{K} B_{K} D_{K} & D_{K} \\
O & D_{K} & O
\end{array}\right)
$$

Since $\operatorname{det} M= \pm\left(\operatorname{det} D_{J}\right)^{2}\left(\operatorname{det} D_{K}\right)^{2} \operatorname{det} B_{J} \neq 0$ the bilinear form is nondegenerate on $H$.

Theorem 16.2 Suppose $A$ is a symmetrisable GCM. Then the Kac-Moody algebra $L(A)$ has a non-degenerate symmetric invariant bilinear form.

Proof. We have $L(A)=\bigoplus_{\alpha \in Q} L_{\alpha}$. For $\alpha=m_{1} \alpha_{1}+\cdots+m_{n} \alpha_{n} \in Q$ we define the height of $\alpha$ by ht $\alpha=m_{1}+\cdots+m_{n}$. Then

$$
L(A)=\bigoplus_{i \in \mathbb{Z}} L_{i}
$$

where $L_{i}$ is the direct sum of all $L_{\alpha}$ with ht $\alpha=i$. Since $\left[L_{\alpha} L_{\beta}\right] \subset L_{\alpha+\beta}$ we have $\left[L_{i} L_{j}\right] \subset L_{i+j}$. Thus $L(A)$ may be considered in this way as a $\mathbb{Z}$-graded Lie algebra.

We define, for each integer $r \geq 0$,

$$
L(r)=\bigoplus_{-r \leq i \leq r} L_{i}
$$

Then we have

$$
H=L(0) \subset L(1) \subset L(2) \subset \cdots
$$

and $\bigcup_{r \geq 0} L(r)=L(A)$.
We have already defined a symmetric bilinear form on $H=L(0)$. We shall extend this definition to give a symmetric bilinear form on $L(r)$ for $r=$ $1,2,3, \ldots$ thus eventually defining such a form on $L(A)$. We shall define the form on $L(r)$ by induction on $r$, assuming it is already defined on $L(r-1)$.

We begin with the case $r=1$. We have

$$
L(1)=\left(\bigoplus_{i=1}^{n} \mathbb{C} f_{i}\right) \oplus H \oplus\left(\bigoplus_{i=1}^{n} \mathbb{C} e_{i}\right)
$$

We define a bilinear form $\langle$,$\rangle on L(1)$ which is uniquely determined by the following rules:

$$
\begin{aligned}
& \langle,\rangle \text { agrees with the form already defined on } H \\
& \left\langle L_{i}, L_{j}\right\rangle=0 \text { unless } i+j=0 \\
& \left\langle e_{i}, f_{i}\right\rangle=\left\langle f_{i}, e_{i}\right\rangle=d_{i} \\
& \left\langle e_{i}, f_{j}\right\rangle=\left\langle f_{j}, e_{i}\right\rangle=0 \quad \text { if } i \neq j
\end{aligned}
$$

This bilinear form on $L(1)$ is clearly symmetric. We show

$$
\langle[x y], z\rangle=\langle x,[y z]\rangle
$$

for all $x, y, z \in L(1)$. In showing this we may assume $x \in L_{i}, y \in L_{j}, z \in L_{k}$ for some $i, j, z \in \mathbb{Z}$ with $|i|,|j|,|k| \leq 1$. We may assume $i+j+k=0$ as otherwise both sides of our required equality are zero. The relation is known already when $i, j, k$ are all 0 . Thus we may assume $i, j, k$ are $1,-1,0$ in some order. There are six possible orders, but it is only necessary to check three of them as the other three follow from them. Thus we show

$$
\begin{aligned}
\left\langle\left[e_{i} h\right], f_{j}\right\rangle & =\left\langle e_{i},\left[h f_{j}\right]\right\rangle \\
\left\langle\left[h e_{i}\right], f_{j}\right\rangle & =\left\langle h,\left[e_{i} f_{j}\right]\right\rangle \\
\left\langle\left[h f_{j}\right], e_{i}\right\rangle & =\left\langle h,\left[f_{j} e_{i}\right]\right\rangle
\end{aligned}
$$

for $h \in H$. Both sides are zero in these relations if $i \neq j$. If $i=j$ the relations are valid because

$$
\left\langle h_{i}, h\right\rangle=d_{i} \alpha_{i}(h) \quad \text { for all } h \in H
$$

This follows from the definition of the form $\langle$,$\rangle on H$. Thus we have

$$
\langle[x y], z\rangle=\langle x,[y z]\rangle \quad \text { for all } x, y, z \in L(1)
$$

Now suppose inductively that a symmetric bilinear form has already been defined on $L(r-1)$ and satisfies:

$$
\begin{aligned}
& \left\langle L_{i}, L_{j}\right\rangle=0 \quad \text { unless } i+j=0 \text { for }|i|,|j| \leq r-1 \\
& \langle[x y], z\rangle=\langle x,[y z]\rangle \quad \begin{array}{l}
\text { for all } x \in L_{i}, y \in L_{j}, z \in L_{k} \text { with }|i|,|j|,|k| \leq r-1 \\
\\
\text { and } i+j+k=0 .
\end{array}
\end{aligned}
$$

We shall show this form can be extended to one on $L(r)$ with analogous properties. We extend the form to $L(r)$ by defining

$$
\left\langle L_{i}, L_{j}\right\rangle=0 \quad \text { unless } i+j=0 \text { for }|i|,|j| \leq r .
$$

We must also define $\langle x, y\rangle=\langle y, x\rangle$ for $x \in L_{r}, y \in L_{-r}$. We assume $r \geq 2$.
Now we have $L(A)=N^{-} \oplus H \oplus N$ with $H=L_{0}, N^{-}=\bigoplus_{i<0} L_{i}, N=$ $\bigoplus_{i>0} L_{i}$. The algebra $N^{-}$is generated by $f_{1}, \ldots, f_{n}$, thus each element of $N^{-}$can be written as a Lie word in $f_{1}, \ldots, f_{n}$, so is a linear combination of Lie monomials in $f_{1}, \ldots, f_{n}$. An element of $L_{-r}$ is a linear combination of Lie monomials in $f_{1}, \ldots, f_{n}$ such that the number of factors in each Lie monomial is $r$. If $r \geq 2$ each Lie monomial is the Lie product of Lie monomials of degree $s, t$ say with $s+t=r$. It follows that each element $y \in L_{-r}$ can be written in the form

$$
y=\sum_{j}\left[c_{j} d_{j}\right]
$$

where $c_{j} \in L_{-u_{j}}, d_{j} \in L_{-v_{j}}$ with $u_{j}>0, v_{j}>0$ and $u_{j}+v_{j}=r$. The expression of $y$ in this form need not be unique.

Given $x \in L_{r}, y \in L_{-r}$ we write $y=\sum_{j}\left[c_{j} d_{j}\right]$ as above and wish to define

$$
\langle x, y\rangle=\sum_{j}\left\langle\left[x c_{j}\right], d_{j}\right\rangle .
$$

The right-hand side is known since $\left[x c_{j}\right]$ and $d_{j}$ lie in $L(r-1)$, so if there is a form of the required type on $L(r)$ it must satisfy the above relation in order to be invariant. However, the right-hand side appears to depend on the particular expression $y=\sum_{j}\left[c_{j} d_{j}\right]$ for $y$ which need not be unique. We must therefore show that the right-hand side remains the same if a different such expression for $y$ is chosen.

In a similar way we can write $x \in L_{r}$ in the form

$$
x=\sum_{i}\left[a_{i} b_{i}\right]
$$

where $a_{i} \in L_{s_{i}}, b_{i} \in L_{t_{i}}$ and $s_{i}>0, t_{i}>0$ with $s_{i}+t_{i}=r$. We shall show

$$
\sum_{i}\left\langle a_{i},\left[b_{i} y\right]\right\rangle=\sum_{j}\left\langle\left[x c_{j}\right], d_{j}\right\rangle .
$$

This will imply that the right-hand side is independent of the given expression for $y$, and also that the left-hand side is independent of the given expression for $x$. In fact it is sufficient to show

$$
\left\langle a_{i},\left[b_{i}\left[c_{j} d_{j}\right]\right]\right\rangle=\left\langle\left[\left[a_{i} b_{i}\right] c_{j}\right], d_{j}\right\rangle .
$$

Now

$$
\begin{aligned}
\left\langle\left[\left[a_{i} b_{i}\right] c_{j}\right], d_{j}\right\rangle & =\left\langle\left[\left[a_{i} c_{j}\right] b_{i}\right], d_{j}\right\rangle-\left\langle\left[\left[b_{i} c_{j}\right] a_{i}\right], d_{j}\right\rangle \\
& =\left\langle\left[a_{i} c_{j}\right],\left[b_{i} d_{j}\right]\right\rangle-\left\langle\left[b_{i} c_{j}\right],\left[a_{i} d_{j}\right]\right\rangle \\
& =\left\langle\left[a_{i} c_{j}\right],\left[b_{i} d_{j}\right]\right\rangle-\left\langle\left[a_{i} d_{j}\right],\left[b_{i} c_{j}\right]\right\rangle \\
& =\left\langle a_{i},\left[c_{j}\left[b_{i} d_{j}\right]\right]\right\rangle-\left\langle a_{i},\left[d_{j}\left[b_{i} c_{j}\right]\right]\right\rangle \\
& =\left\langle a_{i},\left[b_{i}\left[c_{j} d_{j}\right]\right]\right\rangle
\end{aligned}
$$

using the invariance of the form on $L(r-1)$. Hence our form $\langle x, y\rangle$ is now well defined on $L(r)$, where it is bilinear and symmetric.

We must now check that

$$
\langle[x y], z\rangle=\langle x,[y z]\rangle
$$

when $x \in L_{i}, y \in L_{j}, z \in L_{k}$ with $|i|,|j|,|k| \leq r$ and $i+j+k=0$. This is known already by induction unless at least one of $|i|,|j|,|k|$ is equal to $r$.

It is impossible for all of $|i|,|j|,|k|$ to be equal to $r$ since $i+j+k=0$. We suppose first that just one of $|i|,|j|,|k|$ is $r$. Then the other two are non-zero. If $|i|=r$ then

$$
\langle x,[y z]\rangle=\langle[x y], z\rangle
$$

by definition of the form on $L(r)$. Similarly if $|k|=r$ this relation also holds by definition. So suppose $|j|=r$. We may assume that $y$ has the form $y=[a b]$ where $a \in L_{s}, b \in L_{t}, s+t=j$ and $0<|s|<|j|, 0<|t|<|j|$. Then

$$
\begin{aligned}
\langle[x y], z\rangle & =\langle[x[a b]], z\rangle \\
& =\langle[[b x] a], z\rangle+\langle[[x a] b], z\rangle \\
& =\langle[b x],[a z]\rangle+\langle[x a],[b z]\rangle \\
& =\langle[x b],[z a]\rangle+\langle[x a],[b z]\rangle \\
& =\langle x,[b[z a]]\rangle+\langle x,[a[b z]]\rangle \\
& =\langle x,[[a b] z]\rangle \\
& =\langle x,[y z]\rangle
\end{aligned}
$$

using the invariance of the form on $L(r-1)$.
Now suppose that two of $|i|,|j|,|k|$ are equal to $r$. Then $i, j, k$ are $r,-r, 0$ in some order. Thus one of $x, y, z$ lies in $H$.

Suppose $x \in H$. We may again assume $y=[a b]$ where $a \in L_{s}, b \in L_{t}, s+t=$ $j, 0<|s|<|j|, 0<|t|<|j|$.

Then

$$
\begin{aligned}
\langle[x y], z\rangle & =\langle[x[a b]], z\rangle=\langle[[x a] b], z\rangle-\langle[[x b] a], z\rangle \\
& =\langle[x a],[b z]\rangle-\langle[x b],[a z]\rangle \quad \text { by definition of }\langle,\rangle \text { on } L(r) \\
& =\langle x,[a[b z]]\rangle-\langle x,[b[a z]]\rangle \quad \text { by invariance on } L(r-1) \\
& =\langle x,[[a b] z]\rangle=\langle x,[y z]\rangle .
\end{aligned}
$$

If $z \in H$ the result also holds by using the symmetry of the form.
Finally suppose $y \in H$. Then we may assume $z=[a b]$ where $a \in L_{s}$, $b \in L_{t}, s+t=k, 0<|s|<|k|, 0<|t|<|k|$. Then

$$
\begin{aligned}
\langle x,[y z]\rangle & =\langle x,[y[a b]]\rangle=\langle x,[a[y b]]\rangle+\langle x,[[y a] b]\rangle \\
& =\langle[x a],[y b]\rangle+\langle[x[y a]], b\rangle \quad \text { by definition of }\langle,\rangle \text { on } L(r) \\
& =\langle[[x a] y], b\rangle+\langle[x[y a]], b\rangle \quad \text { by invariance on } L(r-1) \\
& =\langle[[x y] a], b\rangle \\
& =\langle[x y],[a b]\rangle \quad \text { by definition of }\langle,\rangle \text { on } L(r) \\
& =\langle[x y], z\rangle .
\end{aligned}
$$

We have therefore proved invariance when $x \in L_{i}, y \in L_{j}, z \in L_{k}$ with $|i|,|j|,|k| \leq r$ and $i+j+k=0$. It follows that invariance holds for all $x, y, z \in L(r)$. By induction the form is therefore invariant on $L(A)$.

Thus we have now defined a symmetric invariant bilinear form on $L(A)$. We show it is non-degenerate. Let $I$ be the kernel of $\langle$,$\rangle , i.e. the set of$ $x \in L(A)$ such that $\langle x, y\rangle=0$ for all $y \in L(A)$. Since the form is invariant $I$ is an ideal of $L(A)$. Since by Proposition 16.1 the form is non-degenerate on restriction to $H$ we have $I \cap H=O$. But the Kac-Moody algebra $L(A)$ has no non-zero ideal $I$ with $I \cap H=O$. Hence $I=O$ and the form is non-degenerate on $L(A)$.

Note The proof of this theorem shows that any symmetric invariant bilinear form on $L(A)$ is uniquely determined by its restriction to $H$.

Definition The form constructed in Theorem 16.2 will be called the standard invariant form on $L(A)$.

Corollary 16.3 For each $i \in \mathbb{Z}$ the pairing $L_{i} \times L_{-i} \rightarrow \mathbb{C}$ given by $x, y \rightarrow$ $\langle x, y\rangle$ is non-degenerate.

Proof. Suppose $x \in L_{i}$ satisfies $\langle x, y\rangle=0$ for all $y \in L_{-i}$. Since $\left\langle L_{i}, L_{j}\right\rangle=0$ unless $i+j=0$ we have $\langle x, y\rangle=0$ for all $y \in L(A)$. Hence $x=0$.

Corollary 16.4 $\left\langle L_{\alpha}, L_{\beta}\right\rangle=0$ unless $\alpha+\beta=0$.
Proof. Suppose $\alpha+\beta \neq 0$ and let $x \in L_{\alpha}, y \in L_{\beta}$. Choose $h \in H$ with

$$
(\alpha+\beta)(h) \neq 0
$$

Then

$$
\langle[x h], y\rangle=\langle x,[h y]\rangle
$$

implies

$$
-\alpha(h)\langle x, y\rangle=\beta(h)\langle x, y\rangle
$$

that is

$$
(\alpha+\beta)(h)\langle x, y\rangle=0 .
$$

Hence $\langle x, y\rangle=0$.
Since the form $\langle$,$\rangle is non-degenerate on H$ it determines a bijection $H^{*} \rightarrow H$ given by $\alpha \rightarrow h_{\alpha}^{\prime}$ where

$$
\left\langle h_{\alpha}^{\prime}, h\right\rangle=\alpha(h) \quad \text { for all } h \in H
$$

Corollary 16.5 (i) Suppose $x \in L_{\alpha}, y \in L_{-\alpha}$, Then $[x y]=\langle x, y\rangle h_{\alpha}^{\prime}$.
(ii) The pairing $L_{\alpha} \times L_{-\alpha} \rightarrow \mathbb{C}$ given by $x, y \rightarrow\langle x, y\rangle$ is non-degenerate.
(iii) For each $x \in L_{\alpha}$ with $x \neq 0$ there exists $y \in L_{-\alpha}$ with $[x y] \neq 0$.

Proof. (i) Consider the element $[x y]-\langle x, y\rangle h_{\alpha}^{\prime} \in H$. For all $h \in H$ we have

$$
\begin{aligned}
\left\langle[x y]-\langle x, y\rangle h_{\alpha}^{\prime}, h\right\rangle & =\langle[x y], h\rangle-\langle x, y\rangle\left\langle h_{\alpha}^{\prime}, h\right\rangle \\
& =\langle x,[y h]\rangle-\alpha(h)\langle x, y\rangle \\
& =0 .
\end{aligned}
$$

Since the form is non-degenerate on $H$ we deduce that $[x y]-\langle x, y\rangle h_{\alpha}^{\prime}=0$.
(ii) Since the form is non-degenerate on $L(A)$ and $\left\langle L_{\alpha}, L_{\beta}\right\rangle=0$ unless $\beta=-\alpha$ the pairing $L_{\alpha} \times L_{-\alpha} \rightarrow \mathbb{C}$ must be non-degenerate.
(iii) For each $x \in L_{\alpha}$ with $x \neq 0$ there exists $y \in L_{-\alpha}$ with $\langle x, y\rangle \neq 0$. Hence $[x y] \neq 0$ by (i).

We now consider to what extent a non-degenerate symmetric invariant bilinear form on $L(A)$ is unique. The following proposition deals with this question.

Proposition 16.6 Suppose A is an indecomposable symmetrisable GCM and $\{$,$\} is a non-degenerate symmetric invariant bilinear form on the Kac-Moody$ algebra $L(A)$. Then there exists a non-zero $\xi \in \mathbb{C}$ such that

$$
\{x, y\}=\xi\langle x, y\rangle \text { for all } x, y \in L(A)^{\prime}
$$

Thus such a form is determined on the subalgebra $L(A)^{\prime}$ up to a non-zero constant.

Proof. The argument of Corollary 16.4 shows that $\left\{L_{\alpha}, L_{\beta}\right\}=0$ whenever $\alpha+$ $\beta \neq 0$. In particular we have $\left\{H, L_{\alpha}\right\}=0$ whenever $\alpha \neq 0$. Since $L(A)=H \oplus$ $\sum_{\alpha \neq 0} L_{\alpha}$ it follows that $\{$,$\} is non-degenerate on restriction to H$. The form $\{$, on $L(A)$ is determined by its restriction to $H$ and by the map $L_{\alpha} \times L_{-\alpha} \rightarrow \mathbb{C}$ given by $x, y \rightarrow\{x, y\}$ for each $\alpha \in \Phi$. The argument of Corollary 16.5 shows that, for $x \in L_{\alpha}, y \in L_{-\alpha}$, we have

$$
[x y]=\{x, y\} k_{\alpha}^{\prime}
$$

where $k_{\alpha}^{\prime}$ is the unique element of $H$ satisfying $\left\{k_{\alpha}^{\prime}, h\right\}=\alpha(h)$ for all $h \in H$.
We therefore have

$$
\left[L_{\alpha} L_{-\alpha}\right]=\mathbb{C} h_{\alpha}^{\prime}=\mathbb{C} k_{\alpha}^{\prime}
$$

for each $\alpha \in \Phi$. Thus there exists a non-zero $\xi_{\alpha} \in \mathbb{C}$ with $h_{\alpha}^{\prime}=\xi_{\alpha} k_{\alpha}^{\prime}$. This implies that

$$
\left\{h_{\alpha}^{\prime}, h\right\}=\xi_{\alpha}\left\langle h_{\alpha}^{\prime}, h\right\rangle \quad \text { for all } h \in H
$$

since both sides are equal to $\xi_{\alpha} \alpha(h)$. Let $\alpha_{i}, \alpha_{j}$ be simple roots. Then we have

$$
\left\{h_{\alpha_{i}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\}=\xi_{\alpha_{i}}\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle
$$

and so by the symmetry of the forms

$$
\xi_{\alpha_{i}}\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle=\xi_{\alpha_{j}}\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle
$$

If $A_{i j} \neq 0$ then $\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{j}}^{\prime}\right\rangle \neq 0$ and we have $\xi_{\alpha_{i}}=\xi_{\alpha_{j}}$. If the GCM $A$ is indecomposable this shows that there exists $\xi \neq 0$ in $\mathbb{C}$ such that $\xi_{\alpha_{i}}=\xi$ for all simple roots $\alpha_{i}$. Thus

$$
\left\{h_{\alpha_{i}}^{\prime}, h\right\}=\xi\left\langle h_{\alpha_{i}}^{\prime}, h\right\rangle \quad \text { for all } h \in H
$$

Now for any $\alpha \in \Phi \quad h_{\alpha}^{\prime}$ is a linear combination of the $h_{\alpha_{i}}^{\prime}$. Hence

$$
\left\{h_{\alpha}^{\prime}, h\right\}=\xi\left\langle h_{\alpha}^{\prime}, h\right\rangle \quad \text { for all } h \in H
$$

Thus $\xi_{\alpha}=\xi$ for all $\alpha \in \Phi$. Using the equations

$$
[x y]=\{x, y\} k_{\alpha}^{\prime}=\langle x, y\rangle h_{\alpha}^{\prime}
$$

for $x \in L_{\alpha}, y \in L_{-\alpha}$ we deduce that

$$
\{x, y\}=\xi\langle x, y\rangle \quad \text { for } x \in L_{\alpha}, y \in L_{-\alpha} .
$$

Now $L(A)^{\prime}$ was defined as the subalgebra of $L(A)$ generated by $e_{1}, \ldots, e_{n}$, $f_{1}, \ldots, f_{n}$. We recall from Proposition 14.21 that $L(A)^{\prime}=[L(A) L(A)]$. It follows that $L(A)^{\prime} \cap H$ is generated by $\left[L_{\alpha} L_{-\alpha}\right]=\mathbb{C} h_{\alpha}^{\prime}$ for all $\alpha \in \Phi$. It follows that

$$
\begin{aligned}
& \left\{h^{\prime}, h\right\}=\xi\left\langle h^{\prime}, h\right\rangle \quad \text { for all } h, h^{\prime} \in L(A)^{\prime} \cap H \\
& \{x, y\}=\xi\langle x, y\rangle \quad \text { for all } x \in L_{\alpha}, y \in L_{-\alpha}
\end{aligned}
$$

But $L(A)^{\prime}=\left(L(A)^{\prime} \cap H\right) \oplus \sum_{\alpha \neq 0} L_{\alpha}$ also by Proposition 14.21. Thus we see that

$$
\{x, y\}=\xi\langle x, y\rangle \quad \text { for all } x, y \in L(A)^{\prime} .
$$

Corollary 16.7 $L(A)^{\prime} \cap H$ is the subspace of $H$ spanned by $h_{1}, \ldots, h_{n}$.
Proof. We saw in the proof of Proposition 16.6 that $L(A)^{\prime} \cap H$ is the subspace generated by the elements $h_{\alpha}^{\prime}$ for all $\alpha \in \Phi$. Each $h_{\alpha}^{\prime}$ is a linear combination of $h_{1}, \ldots, h_{n}$ and so the result follows.

Corollary 16.8 Any non-degenerate symmetric invariant bilinear form on a finite dimensional simple Lie algebra is a constant multiple of the Killing form.

Proof. Since $L(A)$ is simple we have $L(A)^{\prime}=[L(A) L(A)]=L(A)$. Thus the given form is a constant multiple of the Killing form on the whole of $L(A)$.

Important comment on notation. In the case when $L(A)$ has finite type the standard invariant form is not the same as the Killing form. It is a constant multiple of the Killing form.

In our development of the theory of finite dimensional simple Lie algebras we have used the notation $\langle$,$\rangle to denote the Killing form. In the theory of Kac-$ Moody algebras the Killing form does not exist in general, but the standard invariant form exists whenever the Kac-Moody algebra is symmetrisable. In the subsequent development the notation $\langle$,$\rangle will denote the standard$
invariant form of a symmetrisable Kac-Moody algebra. This will be so even in the case of finite dimensional simple Lie algebras, i.e. $\langle$,$\rangle will subsequently$ denote the standard invariant form rather than the Killing form.

### 16.2 The Weyl group of a Kac-Moody algebra

Lemma 16.9 Let $x \in L(A)$ and $J$ be the ideal of $L(A)$ generated by $x$. Then $J=\mathfrak{U}(L(A)) x$.

Proof. The adjoint representation of $L(A)$ gives a Lie algebra homomorphism $L(A) \rightarrow$ [End $L(A)]$. By Proposition 9.3 there is an associative algebra homomorphism

$$
\mathfrak{U}(L(A)) \rightarrow \text { End } L(A)
$$

A subspace $K$ of $L(A)$ satisfies $[L(A), K] \subset K$ if and only if $\mathfrak{H}(L(A)) K \subset$ $K$. Now we have $[L(A), J] \subset J$. Hence $\mathfrak{H}(L(A)) J \subset J$. Since $x \in J$ we have $\mathfrak{H}(L(A)) x \subset J$.

On the other hand $\mathfrak{H}(L(A))(\mathfrak{H}(L(A)) x)=\mathfrak{U}(L(A)) x$, thus

$$
[L(A), \mathfrak{U}(L(A)) x] \subset \mathfrak{H}(L(A)) x
$$

Hence $\mathfrak{H}(L(A)) x$ is an ideal of $L(A)$ containing $x$. Hence $\mathfrak{H}(L(A)) x \supset J$. Thus we must have equality.

Proposition 16.10 In $L(A)$ we have, for $i \neq j,\left(\operatorname{ad} e_{i}\right)^{1-A_{i j}} e_{j}=0$ and $\left(\operatorname{ad} f_{i}\right)^{1-A_{i j}} f_{j}=0$.

Proof. We shall show $\left(\operatorname{ad} f_{i}\right)^{1-A_{i j}} f_{j}=0$. The other relation holds similarly.
Let $x=\left(\operatorname{ad} f_{i}\right)^{1-A_{i j}} f_{j} \in N^{-}$. We shall show $\left[e_{k}, x\right]=0$ for all $k=1, \ldots, n$. Suppose this is so. Then the set of all $y \in L(A)$ with $[y x]=0$ is a subalgebra containing $e_{1}, \ldots, e_{n}$, so contains $N$. Thus $[N, x]=O$ and so $\mathfrak{H}(N) x=$ $\mathbb{C} x$. Since $L(A)=N^{-} \oplus H \oplus N$ we have $\mathfrak{H}(L(A))=\mathfrak{H}\left(N^{-}\right) \mathfrak{H}(H) \mathfrak{H}(N)$ by the PBW basis theorem. Hence

$$
\mathfrak{U}(L(A)) x=\mathfrak{U}\left(N^{-}\right) \mathfrak{U}(H) \mathfrak{U}(N) x=\mathfrak{U}\left(N^{-}\right) \mathfrak{U}(H) x .
$$

Since $\left[H, N^{-}\right] \subset N^{-}$we have $\mathfrak{H}(H) N^{-} \subset N^{-}$and $\mathfrak{H}(H) x \subset N^{-}$. Thus

$$
\mathfrak{H}(L(A)) x \subset \mathfrak{H}\left(N^{-}\right) N^{-} \subset N^{-} .
$$

Let $J=\mathfrak{l}(L(A)) x$. This is the ideal of $L(A)$ generated by $x$, by Lemma 16.9. We have $J \subset N^{-}$so $J \cap H=O$. This implies $J=O$ by definition of $L(A)$. Thus $x=0$.

Thus in order to obtain the required result $x=0$ it is sufficient to show

$$
\left[e_{k},\left(\operatorname{ad} f_{i}\right)^{1-A_{i j}} f_{j}\right]=0 \quad \text { when } i \neq j
$$

We first suppose $k \neq i$ and $k \neq j$. Then

$$
\begin{aligned}
{\left[e_{k},\left(\operatorname{ad} f_{i}\right)^{t} f_{j}\right] } & =\left[e_{k},\left[f_{i},\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}\right]\right] \\
& =\left[\left[e_{k}, f_{i}\right],\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}\right]+\left[f_{i},\left[e_{k},\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}\right]\right] \\
& =\left[f_{i},\left[e_{k},\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}\right]\right]
\end{aligned}
$$

Repeating we obtain, for each $t$,

$$
\left[e_{k},\left(\operatorname{ad} f_{i}\right)^{t} f_{j}\right]=\left(\operatorname{ad} f_{i}\right)^{t}\left[e_{k}, f_{j}\right]=0
$$

Next suppose $k=j$. Then

$$
\left[e_{j},\left(\operatorname{ad} f_{i}\right)^{t} f_{j}\right]=\left(\operatorname{ad} f_{i}\right)^{t}\left[e_{j}, f_{j}\right]=\left(\operatorname{ad} f_{i}\right)^{t} h_{j}
$$

as above. If $1-A_{i j} \geq 2$ then this shows that $\left[e_{j},\left(\operatorname{ad} f_{i}\right)^{1-A_{i j}} f_{j}\right]=0$. If $1-A_{i j}=1$ then $A_{i j}=0$ so $\left[f_{i}, h_{j}\right]=A_{j i} f_{i}=0$. Thus $\left[e_{j},\left(\operatorname{ad} f_{i}\right)^{1-A_{i j}} f_{j}\right]=0$ in this case also.

Finally we suppose $k=i$. Then

$$
\begin{aligned}
{\left[e_{i},\left(\operatorname{ad} f_{i}\right)^{t} f_{j}\right] } & =\left[e_{i},\left[f_{i},\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}\right]\right] \\
& =\left[\left[e_{i}, f_{i}\right],\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}\right]+\left[f_{i},\left[e_{i},\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}\right]\right] \\
& =\left[h_{i},\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}\right]+\left[f_{i},\left[e_{i},\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}\right]\right] \\
& =-\left((t-1) \alpha_{i}+\alpha_{j}\right)\left(h_{i}\right)\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}+\left[f_{i},\left[e_{i},\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}\right]\right] \\
& =\left(-2(t-1)-A_{i j}\right)\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}+\left[f_{i},\left[e_{i},\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}\right]\right]
\end{aligned}
$$

Repeating, we obtain

$$
\begin{aligned}
& \left(-2(t-1)-A_{i j}\right)\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}+\left(-2(t-2)-A_{i j}\right)\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}+\cdots \\
& \quad+\left(-A_{i j}\right)\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}=-t\left(t-1+A_{i j}\right)\left(\operatorname{ad} f_{i}\right)^{t-1} f_{j}
\end{aligned}
$$

We now put $t=1-A_{i j}$. Then we have

$$
\left[e_{i},\left(\operatorname{ad} f_{i}\right)^{1-A_{i j}} f_{j}\right]=0
$$

This completes the proof in all cases.

Using this result of Proposition 16.10 we may deduce, as in Proposition 7.17 , that the maps ad $e_{i}$ and ad $f_{i}$ are locally nilpotent. Then the proof of Proposition 3.4 shows that expad $e_{i}$ and $\exp \operatorname{ad}\left(-f_{i}\right)$ are automorphisms of $L(A)$. Let

$$
n_{i}=\exp \operatorname{ad} e_{i} \cdot \exp \operatorname{ad}\left(-f_{i}\right) \cdot \exp \text { ad } e_{i} \in \text { Aut } L(A)
$$

Proposition $16.11 n_{i}(H)=H$. For $x \in H$ we have

$$
n_{i}(x)=x-\alpha_{i}(x) h_{i}
$$

Proof. Let $x \in H$. Then

$$
\begin{aligned}
& \exp \text { ad } e_{i} \cdot x=\left(1+\operatorname{ad} e_{i}\right) x=x+\left[e_{i} x\right]=x-\alpha_{i}(x) e_{i} \\
& \exp \operatorname{ad}\left(-f_{i}\right) \cdot\left(x-\alpha_{i}(x) e_{i}\right)=\left(1-\operatorname{ad} f_{i}+\frac{\left(\operatorname{ad} f_{i}\right)^{2}}{2}\right)\left(x-\alpha_{i}(x) e_{i}\right) \\
& =x-\alpha_{i}(x) e_{i}-\left[f_{i} x\right]+\alpha_{i}(x)\left[f_{i} e_{i}\right]+\frac{1}{2} \text { ad } f_{i}\left(\left[f_{i} x\right]+\alpha_{i}(x) h_{i}\right) \\
& =x-\alpha_{i}(x) e_{i}-\alpha_{i}(x) f_{i}-\alpha_{i}(x) h_{i}+\frac{1}{2} \alpha_{i}(x) \cdot 2 f_{i} \\
& =x-\alpha_{i}(x) e_{i}-\alpha_{i}(x) h_{i} \\
& \exp \text { ad } e_{i}\left(x-\alpha_{i}(x) e_{i}-\alpha_{i}(x) h_{i}\right)=\left(1+\operatorname{ad} e_{i}\right)\left(x-\alpha_{i}(x) e_{i}-\alpha_{i}(x) h_{i}\right) \\
& =x-\alpha_{i}(x) e_{i}-\alpha_{i}(x) h_{i}+\left[e_{i} x\right]-\alpha_{i}(x)\left[e_{i} h_{i}\right] \\
& =x-\alpha_{i}(x) e_{i}-\alpha_{i}(x) h_{i}-\alpha_{i}(x) e_{i}+2 \alpha_{i}(x) e_{i} \\
& =x-\alpha_{i}(x) h_{i} .
\end{aligned}
$$

This gives the required result.

Proposition 16.12 The map $s_{i}: H \rightarrow H$ induced by $n_{i}$ satisfies $s_{i}^{2}=1$, $s_{i}\left(h_{i}\right)=-h_{i}, \quad s_{i}(x)=x$ when $\left\langle h_{i}, x\right\rangle=0$.

Proof. This follows from $s_{i}(x)=x-\alpha_{i}(x) h_{i}$ together with $\alpha_{i}\left(h_{i}\right)=2$ and $\left\langle h_{i}, x\right\rangle=d_{i} \alpha_{i}(x)$.

The maps $s_{i}: H \rightarrow H$ are called fundamental reflections. The group $W$ of non-singular linear transformations of $H$ generated by $s_{1}, \ldots, s_{n}$ is called the Weyl group $W$ of $L(A)$.

Proposition 16.13 The bilinear form $\langle$,$\rangle on H$ is invariant under $W$.

## Proof. Let $x, y \in H$. Then

$$
\begin{aligned}
\left\langle s_{i} x, s_{i} y\right\rangle & =\left\langle x-\alpha_{i}(x) h_{i}, y-\alpha_{i}(y) h_{i}\right\rangle \\
& =\langle x, y\rangle-\alpha_{i}(x)\left\langle h_{i}, y\right\rangle-\alpha_{i}(y)\left\langle x, h_{i}\right\rangle+\alpha_{i}(x) \alpha_{i}(y)\left\langle h_{i}, h_{i}\right\rangle \\
& =\langle x, y\rangle-\alpha_{i}(x) d_{i} \alpha_{i}(y)-\alpha_{i}(y) d_{i} \alpha_{i}(x)+\alpha_{i}(x) \alpha_{i}(y) \cdot 2 d_{i} \\
& =\langle x, y\rangle .
\end{aligned}
$$

We may also define an action of $W$ on $H^{*}$ by

$$
(w \lambda) x=\lambda\left(w^{-1} x\right) \quad \text { for } w \in W, \lambda \in H^{*}, x \in H
$$

This action is compatible with the isomorphism $H^{*} \rightarrow H$ given by $\lambda \rightarrow h_{\lambda}^{\prime}$ where $\left\langle h_{\lambda}^{\prime}, x\right\rangle=\lambda(x)$ for all $x \in H$. For suppose $w(\lambda)=\mu$ for $\lambda, \mu \in H^{*}$. Then

$$
\begin{aligned}
\left\langle w\left(h_{\lambda}^{\prime}\right), x\right\rangle & =\left\langle h_{\lambda}^{\prime}, w^{-1}(x)\right\rangle=\lambda\left(w^{-1}(x)\right) \\
& =(w \lambda) x=\mu(x)=\left\langle h_{\mu}^{\prime}, x\right\rangle
\end{aligned}
$$

for all $x \in H$. Thus $w\left(h_{\lambda}^{\prime}\right)=h_{\mu}^{\prime}$.

Proposition 16.14 The action of $s_{i}$ on $H^{*}$ is given by

$$
s_{i}(\lambda)=\lambda-\lambda\left(h_{i}\right) \alpha_{i}
$$

Proof. Let $x \in H$. Then

$$
\begin{aligned}
\left(s_{i} \lambda\right) x & =\lambda\left(s_{i}^{-1} x\right)=\lambda\left(s_{i} x\right)=\lambda\left(x-\alpha_{i}(x) h_{i}\right) \\
& =\lambda(x)-\lambda\left(h_{i}\right) \alpha_{i}(x)=\left(\lambda-\lambda\left(h_{i}\right) \alpha_{i}\right) x
\end{aligned}
$$

In fact the Weyl group acts on the root system $\Phi$ of $L(A)$.

Proposition 16.15 If $\alpha \in \Phi, w \in W$ then $w(\alpha) \in \Phi$. Moreover $\operatorname{dim} L_{\alpha}=$ $\operatorname{dim} L_{w(\alpha)}$.

Proof. The proof of Proposition 7.21 also applies in our present situation.

We shall now determine the order of the product $s_{i} s_{j}$ of two distinct fundamental reflections.

Theorem 16.16 Suppose $i \neq j$. Then the order of $s_{i} s_{j} \in W$ is:

| 2 | if | $A_{i j} A_{j i}=0$ |
| :--- | :--- | :--- |
| 3 | if | $A_{i j} A_{j i}=1$ |
| 4 | if | $A_{i j} A_{j i}=2$ |
| 6 | if | $A_{i j} A_{j i}=3$ |
| $\infty$ | if | $A_{i j} A_{j i} \geq 4$. |

Proof. The Weyl group $W$ acts faithfully on $H^{*}$. Let $K$ be the 2-dimensional subspace of $H^{*}$ given by $K=\mathbb{C} \alpha_{i}+\mathbb{C} \alpha_{j}$. We have

$$
\begin{aligned}
& s_{i}\left(\alpha_{i}\right)=-\alpha_{i}, \quad s_{i}\left(\alpha_{j}\right)=\alpha_{j}-A_{i j} \alpha_{i} \\
& s_{j}\left(\alpha_{i}\right)=\alpha_{i}-A_{j i} \alpha_{j}, \quad s_{j}\left(\alpha_{j}\right)=-\alpha_{j} .
\end{aligned}
$$

Thus the subgroup $\left\langle s_{i}, s_{j}\right\rangle$ of $W$ acts on $K$. We obtain a 2-dimensional representation of $\left\langle s_{i}, s_{j}\right\rangle$ given by
$s_{i} \rightarrow\left(\begin{array}{cc}-1 & -A_{i j} \\ 0 & 1\end{array}\right) \quad s_{j} \rightarrow\left(\begin{array}{cc}1 & 0 \\ -A_{j i} & -1\end{array}\right) \quad s_{i} s_{j} \rightarrow\left(\begin{array}{cc}-1+A_{i j} A_{j i} & A_{i j} \\ -A_{j i} & -1\end{array}\right)$.
Consider the order of this $2 \times 2$ matrix representing $s_{i} s_{j}$. Its characteristic polynomial is

$$
\left|\begin{array}{cc}
\lambda+1-A_{i j} A_{j i} & -A_{i j} \\
A_{j i} & \lambda+1
\end{array}\right|=\lambda^{2}+\left(2-A_{i j} A_{j i}\right) \lambda+1
$$

The discriminant of this polynomial is

$$
D=\left(2-A_{i j} A_{j i}\right)^{2}-4=A_{i j} A_{j i}\left(A_{i j} A_{j i}-4\right)
$$

Thus there are two equal eigenvalues if $A_{i j} A_{j i}=0$ or 4 , two distinct complex eigenvalues if $A_{i j} A_{j i}=1,2$ or 3 , and two distinct real eigenvalues if $A_{i j} A_{j i}>4$.

Suppose $A_{i j} A_{j i}=0$. Then $s_{i} s_{j} \rightarrow\left(\begin{array}{cc}-1 & 0 \\ 0 & -1\end{array}\right)$ and the matrix has order 2.
Suppose $A_{i j} A_{j i}=1$. Then the characteristic polynomial is $\lambda^{2}+\lambda+1$ so the eigenvalues are $\omega, \omega^{2}$ where $\omega=\mathrm{e}^{2 \pi \mathrm{i} / 3}$. Thus the matrix is similar to $\left(\begin{array}{cc}\omega & 0 \\ 0 & \omega^{2}\end{array}\right)$ and so has order 3.

Suppose $A_{i j} A_{j i}=2$. The characteristic polynomial is then $\lambda^{2}+1=$ $(\lambda-i)(\lambda+i)$. Thus the matrix is similar to $\left(\begin{array}{cc}i & 0 \\ 0 & -\mathrm{i}\end{array}\right)$ and so has order 4.

Suppose $A_{i j} A_{j i}=3$. The characteristic polynomial is $\lambda^{2}-\lambda+1=$ $(\lambda+\omega)\left(\lambda+\omega^{2}\right)$. Thus the matrix is similar to $\left(\begin{array}{cc}-\omega & 0 \\ 0 & -\omega^{2}\end{array}\right)$ and so has order 6 .

Suppose $A_{i j} A_{j i}=4$. The characteristic polynomial is $\lambda^{2}-2 \lambda+1=(\lambda-1)^{2}$. The eigenvalues are 1,1 , so the matrix is similar to $\left(\begin{array}{ll}1 & \xi \\ 0 & 1\end{array}\right)$ with $\xi \neq 0$. This matrix has infinite order.

Now suppose $A_{i j} A_{j i}>4$. Then the eigenvalues are real and their product is 1. They are also positive and unequal, so have form $\xi, \xi^{-1}$ where $\xi>1$. Thus the matrix is similar to $\left(\begin{array}{cc}\xi & 0 \\ 0 & \xi^{-1}\end{array}\right)$ so has infinite order.

We have so far considered the action of $s_{i} s_{j}$ on the 2-dimensional subspace $K$ of $H^{*}$. We now consider the action of $s_{i} s_{j}$ on the whole of $H^{*}$. Let

$$
K^{\prime}=\left\{\lambda \in H^{*} ; \lambda\left(h_{i}\right)=0, \quad \lambda\left(h_{j}\right)=0\right\}
$$

Then $\operatorname{dim} K^{\prime}=\operatorname{dim} H^{*}-2$. Let $\lambda \in K \cap K^{\prime}$. Then $\lambda=\xi \alpha_{i}+\eta \alpha_{j}$ and

$$
\begin{aligned}
& \lambda\left(h_{i}\right)=2 \xi+\eta A_{i j}=0 \\
& \lambda\left(h_{j}\right)=\xi A_{j i}+2 \eta=0
\end{aligned}
$$

Now $\left|\begin{array}{cc}2 & A_{i j} \\ A_{j i} & 2\end{array}\right|=4-A_{i j} A_{j i}$. Thus if $A_{i j} A_{j i} \neq 4$ we have $\xi=0, \eta=0$ so $K \cap K^{\prime}=$ $O$. Then $H=K \oplus K^{\prime}$. Now $s_{i} s_{j}$ acts trivially on $K^{\prime}$ since, for $\lambda \in K^{\prime}$, we have

$$
s_{i} s_{j}(\lambda)=s_{i}\left(\lambda-\lambda\left(h_{j}\right) \alpha_{j}\right)=s_{i}(\lambda)=\lambda-\lambda\left(h_{i}\right) \alpha_{i}=\lambda
$$

Thus the order of $s_{i} s_{j}$ on $H^{*}$ is equal to the order of $s_{i} s_{j}$ on $K$ provided $A_{i j} A_{j i} \neq 4$. If $A_{i j} A_{j i}=4$ the order of $s_{i} s_{j}$ on $K$ is infinite, so $s_{i} s_{j}$ has infinite order on $H^{*}$.

We now define $l(w)$ and $n(w)$ for $w \in W$ in the same way as when $L(A)$ is finite dimensional. $l(w)$ is the minimal length of $w$ as a product of generators $s_{1}, \ldots, s_{n}$, and $n(w)$ is the number of $\alpha \in \Phi^{+}$with $w(\alpha) \in \Phi^{-}$. Then the proof of Theorem 5.15 also applies in our present situation and shows that $W$ satisfies the deletion condition. Also the proof of Corollary 5.16 applies in our situation and shows that $l(w)=n(w)$. Finally the proof of Theorem 5.18 applies and shows that $W$ is generated by $s_{1}, \ldots, s_{n}$ as a Coxeter group. Thus we have:

Theorem 16.17 The Weyl group $W$ of the Kac-Moody algebra $L(A)$ is a Coxeter group generated by $s_{1}, \ldots, s_{n}$ with relations

$$
\begin{array}{lll}
s_{i}^{2}=1 \\
\left(s_{i} s_{j}\right)^{2}=1 & \text { if } & A_{i j} A_{j i}=0 \\
\left(s_{i} s_{j}\right)^{3}=1 & \text { if } & A_{i j} A_{j i}=1 \\
\left(s_{i} s_{j}\right)^{4}=1 & \text { if } & A_{i j} A_{j i}=2 \\
\left(s_{i} s_{j}\right)^{6}=1 & \text { if } & A_{i j} A_{j i}=3 .
\end{array}
$$

### 16.3 The roots of a Kac-Moody algebra

Let $A$ be a GCM and $L(A)$ the corresponding Kac-Moody algebra. Then

$$
L(A)=H \oplus \sum_{\alpha \in \Phi} L_{\alpha}
$$

where $\Phi=\left\{\alpha \neq 0 ; L_{\alpha} \neq O\right\}$. $\Phi$ is the set of roots of $L(A)$. We recall that $\Phi=\Phi^{+} \cup \Phi^{-}$where $\Phi^{+}=\Phi \cap Q^{+}$and $\Phi^{-}=\Phi \cap Q^{-}$. These are the positive and negative roots. $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ is a subset of $\Phi^{+}$called the set of fundamental roots. The multiplicity of the root $\alpha$ is defined as $\operatorname{dim} L_{\alpha}$. We know from Proposition 14.19 that the fundamental roots $\alpha_{1}, \ldots, \alpha_{n}$ have multiplicity 1 . We also know from Proposition 16.15 that the Weyl group $W$ acts on $\Phi$ and preserves multiplicities.

Definition $\alpha \in \Phi$ is called a real root if there exist $\alpha_{i} \in \Pi$ and $w \in W$ such that $\alpha=w\left(\alpha_{i}\right)$.
$\alpha \in \Phi$ is called an imaginary root if $\alpha$ is not real.
We note that if $\alpha$ is a real root so is $-\alpha$. For let $\alpha=w\left(\alpha_{i}\right)$. Then $-\alpha=$ $w s_{i}\left(\alpha_{i}\right)$. It follows that if $\alpha$ is an imaginary root so is $-\alpha$.

Proposition 16.18 Let $\alpha$ be a real root. Then $\alpha$ has multiplicity 1. Also, for $k \in \mathbb{Z}, k \alpha$ is a root if and only if $k= \pm 1$.

Proof. Since $\alpha=w\left(\alpha_{i}\right)$ and $\alpha_{i}$ has multiplicity 1, Proposition 16.15 implies that $\alpha$ has multiplicity 1 . We also know from Proposition 14.19 that if $k>1$ then $k \alpha_{i}$ is not a root. Since $k \alpha=w\left(k \alpha_{i}\right), k \alpha$ is also not a root.

We now consider the imaginary roots. Let $\Phi_{\mathrm{im}}^{+}$be the set of positive imaginary roots.

Proposition 16.19 If $\alpha \in \Phi_{\mathrm{Im}}^{+}$and $w \in W$ then $w(\alpha) \in \Phi_{\mathrm{Im}}^{+}$.

Proof. We know that $W$ acts both on $\Phi$ and on the set $\Phi_{\mathrm{Re}}$ of real roots. Hence $W$ acts on the set $\Phi_{\mathrm{Im}}$ of imaginary roots. We must show that an element $w \in W$ cannot change the sign of an imaginary root. Let

$$
\alpha=\sum_{i=1}^{n} k_{i} \alpha_{i} \quad k_{i} \geq 0
$$

Now at least two coefficients $k_{i}$ must be positive. Otherwise $\alpha$ would be a multiple of some $\alpha_{i}$ and hence equal to $\alpha_{i}$. But then $\alpha$ would be real, a contradiction. Now $s_{i}(\alpha)=\alpha-\alpha\left(h_{i}\right) \alpha_{i}$, thus $s_{i}(\alpha)$ contains at least one
fundamental root with positive coefficient. Hence $s_{i}(\alpha) \in \Phi_{\mathrm{Im}}^{+}$. Since $w \in W$ is a product of fundamental reflections $s_{i}$ we have $w(\alpha) \in \Phi_{\mathrm{Im}}^{+}$.

We now introduce the fundamental chamber in the context of Kac-Moody algebras. We recall that in Section 12.3 the fundamental chamber was defined for finite dimensional semisimple Lie algebras. In the present context we begin with a GCM $A$ and take a real minimal realisation $\left(H_{\mathbb{R}}, \Pi, \Pi^{v}\right)$ as in Remark 14.20. We then define the fundamental chamber as

$$
C=\left\{\lambda \in H_{\mathbb{R}}^{*} ; \lambda\left(h_{i}\right)>0 \quad \text { for } i=1, \ldots, n\right\}
$$

Its closure is

$$
\bar{C}=\left\{\lambda \in H_{\mathbb{R}}^{*} ; \lambda\left(h_{i}\right) \geq 0 \quad \text { for } i=1, \ldots, n\right\}
$$

Proposition 16.20 Suppose $\alpha \in \Phi_{\operatorname{Im}}^{+}$. Then there exists $w \in W$ with $w(\alpha)$ $\in-\bar{C}$.

Proof. Consider the set of all elements $w(\alpha)$ for $w \in W$. These are all positive imaginary roots by Proposition 16.19. Let $\beta$ be such a root for which ht $\beta$ is as small as possible. Let $\beta=\sum k_{i} \alpha_{i}$. Then $s_{i}(\beta)=\beta-\beta\left(h_{i}\right) \alpha_{i}$. Since ht $s_{i}(\beta) \geq$ ht $\beta$ we have $\beta\left(h_{i}\right) \leq 0$. This holds for all $i$, thus $\beta \in-\bar{C}$.

Proposition 16.21 Let $\alpha \in \Phi, \alpha=\sum_{i=1}^{n} k_{i} \alpha_{i}$ and $\operatorname{supp} \alpha=\left\{i ; k_{i} \neq 0\right\}$. Then supp $\alpha$ is connected.

Proof. We may assume $\alpha \in \Phi^{+}$. Let supp $\alpha=J \subset\{1, \ldots, n\}$. Suppose if possible that $J$ is disconnected, that is $J=J_{1} \cup J_{2}$ with $J_{1}, J_{2}$ non-empty and $A_{i j}=0$ for all $i \in J_{1}, j \in J_{2}$.

We shall show that $\left[e_{i} e_{j}\right]=0$ for all $i \in J_{1}, j \in J_{2}$. We first show the weaker condition

$$
\left[\left[e_{i} e_{j}\right] f_{k}\right]=0 \quad \text { for } i \in J_{1}, j \in J_{2}, k=1, \ldots, n
$$

We have

$$
\begin{aligned}
& {\left[\left[e_{i} e_{j}\right] f_{k}\right]=\left[e_{i}\left[e_{j} f_{k}\right]\right]+\left[\left[e_{i} f_{k}\right] e_{j}\right]} \\
& \text { If } k \notin\{i, j\} \text { then }\left[e_{i} f_{k}\right]=0 \text { and }\left[e_{j} f_{k}\right]=0 . \\
& \text { If } k=i \text { then }\left[\left[e_{i} e_{j}\right] f_{k}\right]=\left[h_{i} e_{j}\right]=A_{i j} e_{j}=0 . \\
& \text { If } k=j \text { then }\left[\left[e_{i} e_{j}\right] f_{k}\right]=\left[e_{i} h_{j}\right]=-A_{j i} e_{i}=0 .
\end{aligned}
$$

Thus in all cases $\left[\left[e_{i} e_{j}\right] f_{k}\right]=0$ for all $k$.

Write $x=\left[e_{i} e_{j}\right]$. The ideal of $L(A)$ generated by $x$ is $\mathfrak{H}(L(A)) x$ as in Lemma 16.9. Since $L(A)=N \oplus H \oplus N^{-}$we have $\mathfrak{H}(L(A))=$ $\mathfrak{H}(N) \mathfrak{U}(H) \mathfrak{U}\left(N^{-}\right)$. Now $N^{-}$is generated by $f_{1}, \ldots, f_{n}$ and $\left[x f_{i}\right]=0$ for each $i$, thus $\mathfrak{U}\left(N^{-}\right) x=\mathbb{C} x$. Since $[H N] \subset N$ we have $\mathfrak{H}(H) N \subset N$ and so $\mathfrak{U}(H) \mathfrak{U}\left(N^{-}\right) x \subset N$ since $x \in N$. Finally $\mathfrak{H}(N) \mathfrak{H}(H) \mathfrak{U}\left(N^{-}\right) x \subset \mathfrak{H}(N) N \subset N$. Thus $\mathfrak{l}(L(A)) x$ is an ideal of $L(A)$ intersecting $H$ in $O$. By definition of $L(A)$ this ideal must be $O$. In particular we have $x=0$. Thus $\left[e_{i} e_{j}\right]=0$ for all $i \in J_{1}, j \in J_{2}$.

We use this fact to obtain the required contradiction. Since $\alpha \in \Phi^{+}$we have $L_{\alpha} \neq O$ and $L_{\alpha} \subset N$. The elements of $L_{\alpha}$ are Lie words in $e_{1}, \ldots, e_{n}$ of weight $\alpha$, and so are linear combinations of Lie monomials in $e_{1}, \ldots, e_{n}$ of weight $\alpha$. Thus there exists a non-zero Lie monomial $m$ in $e_{1}, \ldots, e_{n}$ of weight $\alpha$. We show that any such Lie monomial must be 0 since it contains factors $e_{i}$ both with $i \in J_{1}$ and with $i \in J_{2}$. We can write $m=\left[m_{1} m_{2}\right]$ where $m_{1}, m_{2}$ are shorter Lie monomials. If either $m_{1}$ or $m_{2}$ involves factors $e_{i}$ both with $i \in J_{1}$ and with $i \in J_{2}$ we have $m_{1}=0$ or $m_{2}=0$ by induction. Otherwise all factors $e_{i}$ of $m_{1}$ have $i \in J_{1}$ and all factors $e_{i}$ of $m_{2}$ have $i \in J_{2}$, or vice versa. But then $\left[m_{1} m_{2}\right]=0$ since $\left[e_{i} e_{j}\right]=0$ for all $i \in J_{1}, j \in J_{2}$. Thus $m=0$ and we have the required contradiction.

In order to understand the imaginary roots it will by Proposition 16.20 be sufficient to understand the positive imaginary roots which lie in $-\bar{C}$, the negative of the closure of the fundamental chamber. Such roots satisfy the conditions:

$$
\alpha \in Q^{+}, \alpha \neq 0, \operatorname{supp} \alpha \text { is connected, } \alpha \in-\bar{C} .
$$

It is a remarkable fact that, conversely, any element $\alpha$ satisfying these conditions is a positive imaginary root. Before being able to prove this we need a lemma.

Lemma 16.22 (i) Suppose $\alpha \in \Phi, \alpha \neq \pm \alpha_{i}$, satisfies $\alpha-\alpha_{i} \notin \Phi$ and $\alpha+\alpha_{i} \notin$ Ф. Then $\alpha\left(h_{i}\right)=0$.
(ii) Suppose $\alpha \in \Phi, \alpha \neq-\alpha_{i}$, satisfies $\alpha+\alpha_{i} \notin \Phi$. Then $\alpha\left(h_{i}\right) \geq 0$.

Proof. (i) Since $\alpha \in \Phi$ we have $L_{\alpha} \neq O$. Let $x \in L_{\alpha}$ with $x \neq 0$. Let

$$
n_{i}=\exp \text { ad } e_{i} \cdot \exp \operatorname{ad}\left(-f_{i}\right) \cdot \exp \text { ad } e_{i} \in \operatorname{Aut} L(A)
$$

We show that $n_{i} x \in L_{s_{i}(\alpha)}$. For $[h x]=\alpha(h) x$ for all $h \in H$, hence $\left[n_{i} h, n_{i} x\right]=$ $\alpha(h) n_{i} x$. Now $n_{i}(H)=H$ by Proposition 16.11 and so

$$
\left[h^{\prime}, n_{i} x\right]=\alpha\left(n_{i}^{-1} h^{\prime}\right) n_{i} x \quad \text { for all } h^{\prime} \in H
$$

We have $n_{i}^{-1} h^{\prime}=s_{i}^{-1} h^{\prime}=s_{i} h^{\prime}$ also by Proposition 16.11. Thus $\left[h^{\prime}, n_{i} x\right]=$ $\alpha\left(s_{i} h^{\prime}\right) n_{i} x=\left(s_{i} \alpha\right)\left(h^{\prime}\right) n_{i} x$ for all $h^{\prime} \in H$. Hence $n_{i} x \in L_{s_{i}(\alpha)}$.

Now ad $e_{i} \cdot x \in L_{\alpha+\alpha_{i}}$ and so ad $e_{i} \cdot x=0$ since $\alpha+\alpha_{i} \notin \Phi$ and $\alpha+\alpha_{i} \neq 0$. Thus expad $e_{i} \cdot x=x$. Also ad $f_{i} \cdot x \in L_{\alpha-\alpha_{i}}$ so ad $f_{i} \cdot x=0$ since $\alpha-\alpha_{i} \notin \Phi$ and $\alpha-\alpha_{i} \neq 0$. Thus expad $\left(-f_{i}\right) \cdot x=x$. Hence $n_{i} x=x$. Since $x \in L_{\alpha}$ and $n_{i} x \in L_{s_{i}(\alpha)}$ we deduce $s_{i}(\alpha)=\alpha$. But $s_{i}(\alpha)=\alpha-\alpha\left(h_{i}\right) \alpha_{i}$ and so $\alpha\left(h_{i}\right)=0$.
(ii) Again let $x$ be a non-zero element of $L_{\alpha}$. As before expad $e_{i} \cdot x=x$ since $\alpha+\alpha_{i} \notin \Phi$ and $\alpha+\alpha_{i} \neq 0$. We have

$$
\operatorname{expad}\left(-f_{i}\right) x=x-\operatorname{ad} f_{i} \cdot x+\frac{\left(\operatorname{ad} f_{i}\right)^{2}}{2!} x-\cdots \pm \frac{\left(\operatorname{ad} f_{i}\right)^{p}}{p!} x
$$

where $\left(\operatorname{ad} f_{i}\right)^{p+1} x=0$, since $\operatorname{ad}\left(-f_{i}\right)$ is locally nilpotent. Thus

$$
n_{i} x=\exp \text { ad } e_{i} \cdot \exp \operatorname{ad}\left(-f_{i}\right) \cdot x=\sum_{t \geq 0} \frac{\left(\operatorname{ad} e_{i}\right)^{t}}{t!}\left(x-\operatorname{ad} f_{i} x+\cdots\right)
$$

Now $\left(\operatorname{ad} e_{i}\right)^{t+1}\left(\operatorname{ad} f_{i}\right)^{t} x=0$ for each positive integer $t$, since $\alpha+\alpha_{i} \notin \Phi$ and $\alpha+\alpha_{i} \neq 0$. Hence $\left(\operatorname{ad} e_{i}\right)^{k}\left(\operatorname{ad} f_{i}\right)^{t} x=0$ for all $k \geq t+1$. It follows by considering the above expression for $n_{i} x$ that

$$
n_{i} x \in L_{\alpha} \oplus L_{\alpha-\alpha_{i}} \oplus \cdots \oplus L_{\alpha-p \alpha_{i}}
$$

However, $n_{i} x \in L_{s_{i}(\alpha)}$ as in (i). Thus $s_{i}(\alpha)=\alpha-\alpha\left(h_{i}\right) \alpha_{i}=\alpha-k \alpha_{i}$ for some $k$ with $0 \leq k \leq p$. Hence $\alpha\left(h_{i}\right) \geq 0$.

We now define

$$
K=\left\{\alpha \in Q^{+}, \alpha \neq 0, \operatorname{supp} \alpha \text { is connected, } \alpha \in-\bar{C}\right\}
$$

Proposition 16.23 $K \subset \Phi_{\mathrm{im}}^{+}$.
Proof. Let $\alpha \in K$. Then $\alpha=\sum_{i=1}^{n} k_{i} \alpha_{i}$ with each $k_{i} \geq 0$ and $k_{i}>0$ for some $i$. Also $\operatorname{supp} \alpha=\left\{i ; k_{i} \neq 0\right\}$.

We define a set $\Psi$ of roots by

$$
\Psi=\left\{\beta \in \Phi^{+} ; \beta=\sum_{i=1}^{n} m_{i} \alpha_{i} \quad \text { with } m_{i} \leq k_{i} \quad \text { for each } i\right\}
$$

$\Psi$ is a finite non-empty set of positive roots, since it contains at least one fundamental root. We choose a root $\beta \in \Psi$ such that ht $\beta$ is as large as possible. We aim to show that $\beta=\alpha$ and hence that $\alpha \in \Phi$. We shall show first that $\operatorname{supp} \beta=\operatorname{supp} \alpha$.

Suppose if possible that $\operatorname{supp} \beta \neq \operatorname{supp} \alpha$. Since $\operatorname{supp} \alpha$ is connected there exist $j \in \operatorname{supp} \alpha-\operatorname{supp} \beta$ and $j^{\prime} \in \operatorname{supp} \beta$ such that $A_{j j^{\prime}} \neq 0$. Let $\beta=\sum_{i=1}^{n} m_{i} \alpha_{i}$. Then $m_{j}=0$. Now $\beta-\alpha_{j} \notin \Phi$ since $m_{j}=0$ and $\beta+\alpha_{j} \notin \Phi$ by the maximality of ht $\beta$. By Lemma 16.22 (i) we have $\beta\left(h_{j}\right)=0$. But

$$
\beta\left(h_{j}\right)=\sum_{i \in \operatorname{supp} \beta} m_{i} \alpha_{i}\left(h_{j}\right)=\sum_{i \in \operatorname{supp} \beta} m_{i} A_{j i}<0
$$

since $m_{i}>0, A_{j i} \leq 0$ and $A_{j j^{\prime}}<0$. This contradicts $\beta\left(h_{j}\right)=0$. Thus $\operatorname{supp} \beta=$ $\operatorname{supp} \alpha$.

Now we have $\alpha=\sum_{i=1}^{n} k_{i} \alpha_{i}, \beta=\sum_{i=1}^{n} m_{i} \alpha_{i}$ with $m_{i} \leq k_{i}$. Let

$$
J=\left\{i \in \operatorname{supp} \alpha ; k_{i}=m_{i}\right\}
$$

We aim to show that $J=\operatorname{supp} \alpha$ and so that $\beta=\alpha$.
Suppose if possible that $J \neq \operatorname{supp} \alpha$. Let $i \in \operatorname{supp} \alpha-J$. Then $m_{i}<k_{i}$. Hence $\beta+\alpha_{i} \notin \Phi$ by the maximality of ht $\beta$. Thus $\beta\left(h_{i}\right) \geq 0$ by Lemma 16.22 (ii).

Let $M$ be a connected component of $\operatorname{supp} \alpha-J$. Then $\beta\left(h_{i}\right) \geq 0$ for all $i \in M$. Let $\beta^{\prime}=\sum_{i \in M} m_{i} \alpha_{i}$. Then

$$
\begin{aligned}
\beta^{\prime}\left(h_{i}\right) & =\beta\left(h_{i}\right)-\sum_{j \in \operatorname{supp} \alpha-M} m_{j} \alpha_{j}\left(h_{i}\right) \\
& =\beta\left(h_{i}\right)-\sum_{j \in \operatorname{supp} \alpha-M} m_{j} A_{i j} .
\end{aligned}
$$

If $i \in M$ then $\beta\left(h_{i}\right) \geq 0, m_{j}>0($ since $\operatorname{supp} \beta=\operatorname{supp} \alpha)$ and $A_{i j} \leq 0$. Hence $\beta^{\prime}\left(h_{i}\right) \geq 0$. Since $\operatorname{supp} \alpha$ is connected there exists $i^{\prime} \in M$ and $j^{\prime} \in \operatorname{supp} \alpha-M$ with $A_{i^{\prime} j^{\prime}} \neq 0$. Then

$$
\beta^{\prime}\left(h_{i^{\prime}}\right)=\beta\left(h_{i^{\prime}}\right)-\sum_{j \in \operatorname{supp} \alpha-M} m_{j} A_{i^{\prime} j}
$$

We have $\beta^{\prime}\left(h_{i^{\prime}}\right) \geq 0$ as before; however, in fact we have $\beta^{\prime}\left(h_{i^{\prime}}\right)>0$. The strict inequality holds since $m_{j^{\prime}}>0$ and $A_{i^{\prime} j^{\prime}}<0$.

Let $A_{M}$ be the principal minor $\left(A_{i j}\right)$ with $i, j \in M$. Let $u$ be the column vector with entries $m_{j}$ for $j \in M$. Since

$$
\beta^{\prime}\left(h_{i}\right)=\sum_{j \in M} A_{i j} m_{j} \quad \text { for } i \in M
$$

we have $u>0, A_{M} u \geq 0, A_{M} u \neq 0$. Now we recall that if the indecomposable GCM $A_{M}$ is of affine type then $A_{M} u \geq 0$ implies $A_{M} u=0$. Also if $A_{M}$ is of indefinite type then $A_{M} u \geq 0$ and $u \geq 0$ imply $u=0$. Thus $A_{M}$ cannot have affine type or indefinite type. Hence $A_{M}$ has finite type.

Now let $\gamma=\sum_{i \in \mathrm{M}}\left(k_{i}-m_{i}\right) \alpha_{i}$. We have $k_{i}-m_{i}>0$ for all $i \in M$. We recall that

$$
\alpha-\beta=\sum_{i \in \operatorname{supp} \alpha-J}\left(k_{i}-m_{i}\right) \alpha_{i}
$$

Thus for $i \in M$ we have

$$
(\alpha-\beta)\left(h_{i}\right)=\sum_{j \in \operatorname{supp} \alpha-J}\left(k_{j}-m_{j}\right) A_{i j}=\sum_{j \in M}\left(k_{j}-m_{j}\right) A_{i j}=\gamma\left(h_{i}\right)
$$

since $M$ is a connected component of supp $\alpha-J$. Thus $\gamma\left(h_{i}\right)=\alpha\left(h_{i}\right)-\beta\left(h_{i}\right)$ for all $i \in M$. Now $\alpha\left(h_{i}\right) \leq 0$ since $\alpha \in K$ and $\beta\left(h_{i}\right) \geq 0$ since $i \in M$. Thus $\gamma\left(h_{i}\right) \leq 0$ for all $i \in M$.

Now let $u$ be the column vector with entries $k_{i}-m_{i}$ for $i \in M$. Then we have $u>0$ and $A_{M} u \leq 0$. Since $A_{M}$ has finite type $A_{M}(-u) \geq 0$ implies $-u>0$ or $-u=0$, that is $u<0$ or $u=0$. This is a contradiction since $u>0$. This contradiction shows that $J=\operatorname{supp} \alpha$ and hence that $\beta=\alpha$. Thus $\alpha \in \Phi$. Since $\alpha \in Q^{+}$we have $\alpha \in \Phi^{+}$. Thus $\alpha \in K$ implies $\alpha \in \Phi^{+}$. Now $\alpha \in K$ implies $2 \alpha \in K$, so $2 \alpha \in \Phi^{+}$. By Proposition 16.18 this implies that $\alpha \in \Phi_{\mathrm{Im}}^{+}$. This completes the proof.

This remarkable proof, due to V. Kac, enables us to determine the set of all positive imaginary roots, and hence the set of all imaginary roots.

Theorem 16.24 The set of positive imaginary roots of $L(A)$ is given by

$$
\Phi_{\mathrm{Im}}^{+}=\cup_{w \in W} w(K)
$$

where

$$
K=\left\{\alpha \in Q^{+} ; \alpha \neq 0, \operatorname{supp} \alpha \text { is connected, } \alpha \in-\bar{C}\right\}
$$

The set of all imaginary roots is $\Phi_{\mathrm{Im}}^{+} \cup\left(-\Phi_{\mathrm{Im}}^{+}\right)$.
Proof. This follows from Propositions 16.19, 16.20, 16.21 and 16.23 .

Corollary 16.25 Let $\alpha \in \Phi_{\mathrm{Im}}^{+}$. Then $k \alpha \in \Phi_{\mathrm{Im}}^{+}$for all positive integers $k$.

Proof. This follows from Theorem 16.24 and the fact that $\alpha \in K$ implies $k \alpha \in K$.

We next consider the real and imaginary roots of $L(A)$ when $A$ is symmetrisable. Then $L(A)$ has an invariant bilinear form $\langle$,$\rangle described in Sec-$ tion 16.1. This form is non-degenerate on restriction to $H$, so determines
an isomorphism $H^{*} \rightarrow H$ under which $\lambda \rightarrow h_{\lambda}^{\prime}$, where $\lambda(x)=\left\langle h_{\lambda}^{\prime}, x\right\rangle$ for all $x \in H$. We can then transfer the bilinear form to $H^{*}$ by defining

$$
\langle\lambda, \mu\rangle=\left\langle h_{\lambda}^{\prime}, h_{\mu}^{\prime}\right\rangle .
$$

In particular we can define $\langle\alpha, \alpha\rangle$ for $\alpha \in \Phi$.

Proposition 16.26 Suppose $A$ is a symmetrisable GCM. Then if $\alpha$ is a real root of $L(A)$ we have $\langle\alpha, \alpha\rangle>0$. If $\alpha$ is an imaginary root then $\langle\alpha, \alpha\rangle \leq 0$.

Proof. The form $\langle$,$\rangle on H$ is $W$-invariant by Proposition 16.13, so the induced form on $H^{*}$ is also $W$-invariant. By definition of the form on $H$ we have

$$
\left\langle h_{i}, x\right\rangle=d_{i} \alpha_{i}(x) \quad \text { for all } x \in H
$$

Hence $d_{i} \alpha_{i} \in H^{*}$ corresponds to $h_{i} \in H$ under our map $H^{*} \rightarrow H$. Thus

$$
\left\langle\alpha_{i}, \alpha_{i}\right\rangle=\frac{1}{d_{i}^{2}}\left\langle h_{i}, h_{i}\right\rangle=\frac{2}{d_{i}} .
$$

In particular $\left\langle\alpha_{i}, \alpha_{i}\right\rangle>0$. Now each real root has form $w\left(\alpha_{i}\right)$ for some $w \in W$ and some $i$. Hence

$$
\left\langle w\left(\alpha_{i}\right), w\left(\alpha_{i}\right)\right\rangle=\left\langle\alpha_{i}, \alpha_{i}\right\rangle>0
$$

Now consider the imaginary roots. Let $\alpha \in K$. Then $\alpha=\sum k_{i} \alpha_{i}$ with each $k_{i} \geq 0$ and $\alpha\left(h_{i}\right) \leq 0$ for each $i$. Thus

$$
\langle\alpha, \alpha\rangle=\sum k_{i}\left\langle\alpha, \alpha_{i}\right\rangle=\sum k_{i} \cdot \frac{1}{d_{i}} \alpha\left(h_{i}\right) \leq 0
$$

since $k_{i} \geq 0, d_{i}>0, \alpha\left(h_{i}\right) \leq 0$. Every positive imaginary root has form $w(\alpha)$ for some $w \in W, \alpha \in K$, thus

$$
\langle w(\alpha), w(\alpha)\rangle=\langle\alpha, \alpha\rangle \leq 0
$$

For the negative imaginary roots we have

$$
\langle-w(\alpha),-w(\alpha)\rangle=\langle w(\alpha), w(\alpha)\rangle \leq 0
$$

We next obtain information about the imaginary roots in the three cases of our trichotomy.

Theorem 16.27 Let A be an indecomposable GCM.
(i) If A has finite type then $L(A)$ has no imaginary roots.
(ii) Suppose A has affine type. Then there exists $u>0$ with $A u=0$. The vector $u$ is determined to within a scalar multiple. Thus there is a unique such $u$ whose entries are positive integers with no common factor. Let $u=\left(a_{1}, \ldots, a_{n}\right)$. Let $\delta=a_{1} \alpha_{1}+\cdots+a_{n} \alpha_{n}$. Then the imaginary roots of $L(A)$ are the elements $k \delta$ for $k \in \mathbb{Z}, k \neq 0$.
(iii) Suppose A has indefinite type. Then there exists $\alpha \in \Phi_{\operatorname{Im}}^{+}$such that $\alpha=$ $\sum_{i=1}^{n} k_{i} \alpha_{i}$ with $k_{i}>0$ and $\alpha\left(h_{i}\right)<0$ for all $i=1, \ldots, n$.

Proof. (i) If $A$ has finite type $L(A)$ is a finite dimensional simple Lie algebra by Theorem 15.19. Thus each root of $L(A)$ is real by Proposition 5.12.
(ii) Suppose $A$ has affine type. We first consider the imaginary roots in $K$. Let $\alpha \in K$ satisfy $\alpha=\sum_{i=1}^{n} k_{i} \alpha_{i}$. Let $v$ be the column vector $\left(k_{1}, \ldots, k_{n}\right)$. Then we have $v \geq 0$ and $A v \leq 0$, since $\alpha\left(h_{i}\right) \leq 0$ for each $i$. But in affine type $A(-v) \geq 0$ implies $A(-v)=0$. Thus $v \neq 0$ and $A v=0$.

We also have $u>0$ and $A u=0$. Since $A$ has corank $1, v$ is a multiple of $u$. Since the coefficients of $u$ have no common factor $v=k u$ for some $k \in \mathbb{Z}$ with $k>0$. Thus $\alpha=k \delta$.

Now every positive imaginary root has form $w(\alpha)$ for some $\alpha \in K$, by Theorem 16.24. We have

$$
s_{i}(\delta)=\delta-\delta\left(h_{i}\right) \alpha_{i}=\delta
$$

since $\delta\left(h_{i}\right)=0$ follows from $A u=0$. It follows that $w(\delta)=\delta$ for each $w \in W$. Thus the only positive imaginary roots are the elements $k \delta$ with $k \in \mathbb{Z}, k>0$. Hence the only imaginary roots are the $k \delta$ with $k \in \mathbb{Z}, k \neq 0$.
(iii) Suppose $A$ has indefinite type. Then there exists $u>0$ with $A u<0$. Suppose $u=\left(k_{1}, \ldots, k_{n}\right)$. Let $\alpha=\sum_{i=1}^{n} k_{i} \alpha_{i}$. Then $\alpha \in K$ and $\alpha\left(h_{i}\right)<0$ for all $i$. Thus $\alpha$ is a positive imaginary root of the required kind.

A significant consequence of the last result is as follows.

Corollary 16.28 If $A$ is an indecomposable GCM of affine or indefinite type then the dimension of $L(A)$ is infinite.

Proof. In both cases $L(A)$ has an imaginary root $\alpha$. Thus it has infinitely many imaginary roots $k \alpha$ for $k \in \mathbb{Z}, k \neq 0$, by Corollary 16.25. Since $L(A)=$ $H \oplus \sum_{\alpha \in \Phi} L_{\alpha}$, $\operatorname{dim} L(A)$ must be infinite.

We next consider which of the imaginary roots of $L(A)$ when $A$ is symmetrisable satisfy $\langle\alpha, \alpha\rangle=0$.

Proposition 16.29 Let $A$ be symmetrisable and $\alpha$ be an imaginary root of $L(A)$. Then $\langle\alpha, \alpha\rangle=0$ if and only if there exists $w \in W$ such that the support of $w(\alpha)$ has a diagram of affine type.

Proof. First suppose $\langle\alpha, \alpha\rangle=0$. We may assume without loss of generality that $\alpha \in \Phi^{+}$. Thus there exists $w \in W$ with $w(\alpha) \in K$, by Theorem 16.24. Let $\beta=w(\alpha)$. Then $\beta\left(h_{i}\right) \leq 0$ for all $i$. Let $J$ be the support of $\beta$ and $\beta=\sum_{i \in J} k_{i} \alpha_{i}$. Then $J$ is connected. Now

$$
\langle\beta, \beta\rangle=\sum_{i \in J} k_{i}\left\langle\beta, \alpha_{i}\right\rangle=\sum_{i \in J} \frac{k_{i}}{d_{i}} \beta\left(h_{i}\right) .
$$

Now $k_{i}>0, d_{i}>0$ and $\beta\left(h_{i}\right) \leq 0$ for all $i \in J$. Since we also have $\langle\beta, \beta\rangle=$ $\langle w(\alpha), w(\alpha)\rangle=\langle\alpha, \alpha\rangle=0$ we deduce that $\beta\left(h_{i}\right)=0$ for all $i \in J$. Hence $\sum_{j \in J} k_{j} \alpha_{j}\left(h_{i}\right)=0$, that is $\sum_{j \in J} A_{i j} k_{j}=0$. Let $u$ be the column vector with entries $k_{j}$ for $j \in J$. Then $u>0$ and $A_{J} u=0$. Since $A_{J}$ is an indecomposable GCM this implies that $A_{J}$ has affine type, by Corollary 15.11.

Conversely suppose $\beta$ is a positive imaginary root whose support $J$ has a diagram of affine type. Then $L_{\beta} \neq O$ and so $L(A)$ contains a non-zero Lie monomial in $e_{1}, \ldots, e_{n}$ of weight $\beta$. The letters $e_{i}$ in this Lie monomial all have $i \in J$. Thus the Lie monomial lies in $L\left(A_{J}\right)$ and so $\beta$ is a root of $L\left(A_{J}\right)$. If $\beta$ were a real root of $L\left(A_{J}\right)$ it would have form $w\left(\alpha_{i}\right)$ for some $w \in W\left(A_{J}\right)$ and $i \in J$, and so $\beta$ would be a real root of $L(A)$. Thus $\beta$ is an imaginary root of $L\left(A_{J}\right)$. Since $L\left(A_{J}\right)$ has affine type, $\beta=k \delta$ where $\delta$ is the element for $L\left(A_{J}\right)$ defined in Theorem 16.27 (ii). Let $\delta=\sum_{i \in J} a_{i} \alpha_{i}$. Then

$$
\langle\delta, \delta\rangle=\sum_{i \in J} a_{i}\left\langle\delta, \alpha_{i}\right\rangle=\sum_{i \in J} \frac{a_{i}}{d_{i}} \delta\left(h_{i}\right)=0
$$

since $\delta\left(h_{i}\right)=0$ for all $i \in J$. Thus $\langle\beta, \beta\rangle=0$ also. Finally if $\alpha$ is any root of $L(A)$ satisfying $w(\alpha)=\beta$ for some $w \in W$, we have $\langle\alpha, \alpha\rangle=0$ also.

## 17

## Kac-Moody algebras of affine type

### 17.1 Properties of the affine Cartan matrix

We now consider the Kac-Moody algebras $L(A)$ where $A$ is a GCM of affine type. Let $A$ be an $n \times n$ matrix of rank $l$. Then we know that $n=l+1$. We shall number the rows and columns of $A$ by the integers $0,1, \ldots, l$. There exists a unique vector $a=\left(a_{0}, a_{1}, \ldots, a_{l}\right)$ whose coordinates are positive integers with no common factor such that

$$
A\left(\begin{array}{c}
a_{0} \\
a_{1} \\
\vdots \\
a_{l}
\end{array}\right)=\left(\begin{array}{c}
0 \\
0 \\
\vdots \\
0
\end{array}\right)
$$

The possible Dynkin diagrams of such matrices $A$ were obtained on the affine list 15.20 . We shall choose the numbering of the vertices in such a way that node 0 is the one in black in the diagram below. We also show in each diagram the integer $a_{i}$ associated to each vertex.
17.1 The integers $a_{0}, a_{1}, \ldots, a_{l}$.

$\tilde{A}_{1}$

$\tilde{A}_{2}$

$\tilde{A}_{3}$




$\tilde{D}_{4}$

$\tilde{D}_{5}$

$\tilde{D}_{6}$



There exists also a unique vector $\left(c_{0}, c_{1}, \ldots, c_{l}\right)$ whose coordinates are positive integers with no common factor such that

$$
\left(c_{0}, c_{1}, \ldots, c_{l}\right) A=(0,0, \ldots, 0) .
$$

In fact the vector $\left(c_{0}, c_{1}, \ldots, c_{l}\right)$ for $A$ is the same as the vector $\left(a_{0}, a_{1}, \ldots, a_{l}\right)$ for the transpose $A^{\mathrm{t}}$. Thus the vector $\left(c_{0}, c_{1}, \ldots, c_{l}\right)$ may also be read off from the diagrams in the list 17.1.

Proposition 17.2 (i) $c_{0}=1$.
(ii) $a_{0}=1$ unless $A$ has type $\tilde{C}_{l}^{\prime}$ or $\tilde{A}_{1}^{\prime}$. In these cases $a_{0}=2$.

Proof. This is clear from 17.2

Let $\left(H, \Pi, \Pi^{\mathrm{v}}\right)$ be a minimal realisation of $A$. Then $\operatorname{dim} H=2 n-l=$ $l+2 . \Pi^{\mathrm{v}}=\left\{h_{0}, h_{1}, \ldots, h_{l}\right\}$ is a linearly independent subset of $H$ and $\Pi=$ $\left\{\alpha_{0}, \alpha_{1}, \ldots, \alpha_{l}\right\}$ is a linearly independent subset of $H^{*}$. These exists an element $d \in H$ such that

$$
\alpha_{0}(d)=1 \quad \alpha_{i}(d)=0 \quad \text { for } i=1, \ldots, l .
$$

$d$ is called a scaling element.

Proposition 17.3 $h_{0}, h_{1}, \ldots, h_{l}, d$ is a basis of $H$.

Proof. We must show that $d$ is not a linear combination of $h_{0}, h_{1}, \ldots, h_{l}$. Suppose if possible that $d=\sum_{i=0}^{l} k_{i} h_{i}$. Then $\alpha_{j}(d)=\sum_{i=0}^{l} k_{i} \alpha_{j}\left(h_{i}\right)=\sum_{i=0}^{l} k_{i} A_{i j}$. Hence

$$
\sum_{i=0}^{l} k_{i}\left(A_{i 0}, \ldots, A_{i l}\right)=(1,0, \ldots, 0) .
$$

In particular, omitting the first column of $A$,

$$
\sum_{i=0}^{l} k_{i}\left(A_{i 1}, \ldots, A_{i l}\right)=(0, \ldots, 0)
$$

However, we also have

$$
\sum_{i=0}^{l} c_{i}\left(A_{i 1}, \ldots, A_{i l}\right)=(0, \ldots, 0)
$$

Since the $(l+1) \times l$ matrix $\left(A_{i j}\right), 0 \leq i \leq l, 1 \leq j \leq l$ has rank $l$, this implies that $\left(k_{0}, \ldots, k_{l}\right)$ is a scalar multiple of $\left(c_{0}, \ldots, c_{l}\right)$. But this would imply that

$$
\left(k_{0}, k_{1}, \ldots, k_{l}\right) A=(0,0, \ldots, 0)
$$

a contradiction.
We now define an element $\gamma \in H^{*}$ determined uniquely by

$$
\gamma\left(h_{0}\right)=1 \quad \gamma\left(h_{i}\right)=0 \quad \text { for } i=1, \ldots, l \quad \gamma(d)=0
$$

Proposition $17.4 \alpha_{0}, \alpha_{1}, \ldots, \alpha_{l}, \gamma$ is a basis of $H^{*}$.

Proof. The $(l+2) \times(l+2)$ matrix obtained by applying these elements of $H^{*}$ to the basis of $H$ in Proposition 17.3 is

$$
\begin{aligned}
& 01 \text {. . . } l l+1
\end{aligned}
$$

where $A^{0}$ is a Cartan matrix of finite type. Thus $\operatorname{det} A^{0} \neq 0$ and so the determinant of the above matrix is also non-zero. Hence $\alpha_{0}, \alpha_{1}, \ldots, \alpha_{l}, \gamma$ must be a basis of $H^{*}$.

We know that any indecomposable GCM of affine type is symmetrisable. We shall now express the affine Cartan matrix $A$ in an explicit way as the product of a diagonal matrix $D$ with positive diagonal entries and a symmetric matrix $B$. The diagonal entries of $D$ are rational, but not necessarily integral.

Proposition 17.5 We have $A=D B$ where $D=\operatorname{diag}\left(d_{0}, d_{1}, \ldots, d_{l}\right)$ and $B$ is symmetric, where $d_{i}=a_{i} / c_{i}$.

Proof. By Theorem 15.17 there exists a diagonal matrix $D$ with positive diagonal entries and a symmetric matrix $B$ such that $A=D B$. Let $c=\left(c_{0}, c_{1}, \ldots, c_{l}\right)$ and $a^{\mathrm{t}}=\left(a_{0}, a_{1}, \ldots, a_{l}\right)$. Then $A a=0$ so $D B a=0$, and hence $B a=0$. Thus $a^{\mathrm{t}} B=0$. Also $c A=0$ so $(c D) B=0$. Since $B$ has corank $1 c D$ must be a scalar multiple of $a^{\mathrm{t}}$. In fact we can choose $D$ so that $c D=a^{\mathrm{t}}$, that is $d_{i}=a_{i} / c_{i} . \quad \square$

Now we have a non-degenerate bilinear form on $H$ defined as in Proposition 16.1. This form satisfies

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=d_{i} d_{j} B_{i j}=a_{j} c_{j}^{-1} A_{i j} \quad \text { for } i, j=0,1, \ldots, l \\
& \left\langle h_{0}, d\right\rangle=d_{0} \alpha_{0}(d)=a_{0} \\
& \left\langle h_{i}, d\right\rangle=0 \quad \text { for } i=1, \ldots, l \\
& \langle d, d\rangle=0
\end{aligned}
$$

This standard invariant form on $H$ defines a bijection $H^{*} \rightarrow H$ given by $\lambda \rightarrow h_{\lambda}^{\prime}$ where $\lambda(x)=\left\langle h_{\lambda}^{\prime}, x\right\rangle$ for all $x \in H$.

Proposition 17.6 Under this bijection between $H$ and $H^{*}, h_{i} \in H$ corresponds to $a_{i} c_{i}^{-1} \alpha_{i} \in H^{*}$ for $i=0,1, \ldots, l$ and $d \in H$ corresponds to $a_{0} \gamma \in H^{*}$.

Proof. For $j=0,1, \ldots, l$ we have

$$
\begin{aligned}
& a_{j} c_{j}^{-1} \alpha_{j}\left(h_{i}\right)=d_{j} A_{i j}=\left\langle h_{j}, h_{i}\right\rangle \quad \text { for } i=0,1, \ldots, l \\
& a_{j} c_{j}^{-1} \alpha_{j}(d)=d_{j} \alpha_{j}(d)=\left\langle h_{j}, d\right\rangle
\end{aligned}
$$

thus $a_{j} c_{j}^{-1} \alpha_{j} \in H^{*}$ corresponds to $h_{j} \in H$. We also have

$$
\begin{aligned}
& a_{0} \gamma\left(h_{i}\right)=\left\langle d, h_{i}\right\rangle \quad \text { for } i=0,1, \ldots, l \\
& a_{0} \gamma(d)=\langle d, d\rangle
\end{aligned}
$$

thus $a_{0} \gamma \in H^{*}$ corresponds to $d \in H$.

We may transfer the standard bilinear form from $H$ to $H^{*}$ using this bijection. The form on $H^{*}$ is then given by

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=a_{0}^{-1} \\
& \left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, l \\
& \langle\gamma, \gamma\rangle=0
\end{aligned}
$$

We note in particular that

$$
A_{i j}=\frac{2\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}
$$

Corollary 17.7 Under the given bijection between $H$ and $H^{*}, h_{i} \in H$ corresponds to $\frac{2 \alpha_{i}}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} \in H^{*}$.

Proof. This follows from Proposition 17.6.
We now define an element $c \in H$ by $c=\sum_{i=0}^{l} c_{i} h_{i}$. Under the bijection $H \rightarrow H^{*} c$ corresponds to $\delta$. For $\delta=\sum_{i=0}^{l} a_{i} \alpha_{i}$ and $h_{i}$ corresponds to $a_{i} c_{i}^{-1} \alpha_{i}$ by Proposition 17.6.

Proposition 17.8 The element $c$ lies in the centre of $L(A)$. In fact the centre is 1-dimensional and consists of all scalar multiples of $c$.

Proof. For each simple root $\alpha_{j}$ we have $\alpha_{j}(c)=\sum_{i=0}^{l} c_{i} \alpha_{j}\left(h_{i}\right)=\sum_{i=0}^{l} c_{i} A_{i j}=0$. It follows that $\alpha(c)=0$ for all $\alpha \in \Phi$. Now $L(A)=H \oplus \sum_{\alpha \in \Phi} L_{\alpha}$. Thus each element of $L(A)$ has form $h+\sum x_{\alpha}$ where $h \in H, x_{\alpha} \in L_{\alpha}$ and finitely many $x_{\alpha}$ are non-zero. Thus

$$
\left[c, h+\sum_{\alpha} x_{\alpha}\right]=\sum_{\alpha} \alpha(c) x_{\alpha}=0 .
$$

Hence $c$ lies in the centre of $L(A)$.
Now let $h+\sum_{\alpha} x_{\alpha}$ be any element of the centre of $L(A)$. Then we have

$$
\left[x, h+\sum_{\alpha} x_{\alpha}\right]=0 \quad \text { for all } x \in H
$$

Thus $\sum_{\alpha} \alpha(x) x_{\alpha}=0$ for all $x \in H$. This implies $\alpha(x) x_{\alpha}=0$ for all $x \in H$. Now for each $\alpha \in \Phi$ there exists $x \in H$ with $\alpha(x) \neq 0$. Hence $x_{\alpha}=0$. This shows that the centre of $L(A)$ lies in $H$.

So let $h \in H$ lie in the centre of $L(A)$. By Proposition 17.3 we have

$$
h=\sum_{i=0}^{l} \xi_{i} h_{i}+\xi d \quad \text { for } \xi_{i}, \xi \in \mathbb{C}
$$

Let $x \in L_{\alpha_{j}}$. Then $[h x]=\alpha_{j}(h) x$. Hence $\alpha_{j}(h)=0$ for each $j=0,1, \ldots, l$. Thus

$$
\sum_{i=0}^{l} \xi_{i} \alpha_{j}\left(h_{i}\right)+\xi \alpha_{j}(d)=0
$$

that is

$$
\sum_{i=0}^{l} \xi_{i} A_{i j}=0 \quad \text { for } j=1, \ldots, l
$$

and

$$
\sum_{i=0}^{l} \xi_{i} A_{i 0}=-\xi
$$

However, we have $\sum_{j=0}^{l} A_{i j} a_{j}=0$, hence $A_{i 0}=-a_{0}^{-1} \sum_{j=1}^{l} A_{i j} a_{j}$. Thus $\sum_{i=0}^{l} \xi_{i} A_{i j}=0$ for $j=1, \ldots, l$ implies $\sum_{i=0}^{l} \xi_{i} A_{i 0}=0$. Thus we deduce that $\xi=0$, and so $h=\sum_{i=0}^{l} \xi_{i} h_{i}$. This in turn gives $\sum_{i=0}^{l} \xi_{i} A_{i j}=0$ for $j=0,1, \ldots, l$. This implies that $\left(\xi_{0}, \xi_{1}, \ldots, \xi_{l}\right)$ is a scalar multiple of $\left(c_{0}, c_{1}, \ldots, c_{l}\right)$ since $A$ is an $(l+1) \times(l+1)$ matrix of rank $l$. Thus $h$ is a multiple of $c$. Thus the centre of $L(A)$ is the 1 -dimensional subspace spanned by $c$.
$c$ is called the canonical central element of $L(A)$.

## Summary

We will find it convenient to summarise in one place the properties of the various elements discussed in this section:
(a) $h_{0}, h_{1}, \ldots, h_{l}, d$ are a basis of $H$.
(b) $c=c_{0} h_{0}+\cdots+c_{l} h_{l}$ is the canonical central element.
(c) $\alpha_{0}, \alpha_{1}, \ldots, \alpha_{l}, \gamma$ are a basis of $H^{*}$.
(d) $\delta=a_{0} \alpha_{0}+\cdots+a_{l} \alpha_{l}$ is the basic imaginary root.
(e) The standard invariant form on $H$ is given by

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle h_{0}, d\right\rangle=a_{0} \\
& \left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, l \\
& \langle d, d\rangle=0 .
\end{aligned}
$$

(f) The standard invariant form on $H^{*}$ is given by

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=a_{0}^{-1} \\
& \left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, l \\
& \langle\gamma, \gamma\rangle=0
\end{aligned}
$$

(g) The action of $H^{*}$ on $H$ is given by

$$
\begin{aligned}
& \alpha_{j}\left(h_{i}\right)=A_{i j} \quad i, j=0,1, \ldots, l \\
& \alpha_{0}(d)=1 \\
& \alpha_{i}(d)=0 \quad i=1, \ldots, l \\
& \gamma\left(h_{0}\right)=1 \\
& \gamma\left(h_{i}\right)=0 \quad i=1, \ldots, l \\
& \gamma(d)=0 .
\end{aligned}
$$

(h) The properties of the central element $c$.

$$
\begin{array}{ll}
\left\langle h_{i}, c\right\rangle=0 & i=0,1, \ldots, l \\
\langle d, c\rangle=a_{0} & \\
\langle c, c\rangle=0 & \\
\alpha_{j}(c)=0 & j=0,1, \ldots, l \\
\gamma(c)=1
\end{array}
$$

(i) The properties of the imaginary root $\delta$.

$$
\begin{aligned}
& \left\langle\alpha_{j}, \delta\right\rangle=0 \quad j=0,1, \ldots, l \\
& \langle\gamma, \delta\rangle=1 \\
& \langle\delta, \delta\rangle=0 \\
& \delta\left(h_{i}\right)=0 \quad i=0,1, \ldots, l \\
& \delta(d)=a_{0} \\
& \delta(c)=0
\end{aligned}
$$

(j) Properties of the standard bijection $H \rightarrow H^{*}$.

$$
\begin{aligned}
& h_{i} \rightarrow a_{i} c_{i}^{-1} \alpha_{i} \quad i=0,1, \ldots, l \\
& d \rightarrow a_{0} \gamma \\
& c \rightarrow \delta
\end{aligned}
$$

### 17.2 The roots of an affine Kac-Moody algebra

Let $A^{0}$ be the matrix obtained from the affine Cartan matrix $A$ by removing the row and the column 0 . Then $A^{0}$ is an $l \times l$ Cartan matrix of finite type. By list 17.1 we see that $A^{0}$ is given in each case by the following list.

## The underlying Cartan matrix $A^{0}$

| $\boldsymbol{A}$ |  | $\boldsymbol{A}^{0}$ |
| :--- | :--- | :--- |
| $\tilde{A}_{l}$ | $l \geq 1$ | $A_{l}$ |
| $\tilde{A}_{1}^{\prime}$ |  | $A_{1}$ |
| $\tilde{B}_{l}$ | $l \geq 3$ | $B_{l}$ |
| $\tilde{B}_{l}^{\mathrm{t}}$ | $l \geq 3$ | $C_{l}$ |
| $\tilde{C}_{l}$ | $l \geq 2$ | $C_{l}$ |
| $\tilde{C}_{l}^{\mathrm{t}}$ | $l \geq 2$ | $B_{l}$ |
| $\tilde{C}_{l}^{\prime}$ | $l \geq 2$ | $C_{l}$ |
| $\tilde{D}_{l}$ | $l \geq 4$ | $D_{l}$ |
| $\tilde{E}_{l}$ | $l=6,7,8$ | $E_{l}$ |
| $\tilde{F}_{4}$ |  | $F_{4}$ |
| $\tilde{F}_{4}^{\mathrm{t}}$ |  | $F_{4}$ |
| $\tilde{G}_{2}$ |  | $G_{2}$ |
| $\tilde{G}_{2}^{\mathrm{t}}$ |  | $G_{2}$ |

Let $\Phi^{0}$ be the set of roots of the finite dimensional Lie algebra $L\left(A^{0}\right) . \Phi^{0}$ has a fundamental system $\Pi^{0}=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$. Let $W^{0}$ be the Weyl group of $\Phi^{0}$. Then $W^{0}$ is generated by the fundamental reflections $s_{1}, \ldots, s_{l}$.

Now we know that the imaginary roots of $L(A)$ are the elements $k \delta$ with $k \in \mathbb{Z}$ and $k \neq 0$, by Theorem 16.27 (ii). (However, we do not yet know the multiplicities of these roots.) Thus we shall now consider the real roots of $L(A)$. These have the form $w\left(\alpha_{i}\right)$ for some $w \in W$ and $i=0,1, \ldots, l$. We consider the squared lengths $\langle\alpha, \alpha\rangle$ of the roots $\alpha \in \Phi_{\mathrm{Re}}$. Since $\left\langle w\left(\alpha_{i}\right), w\left(\alpha_{i}\right)\right\rangle=$ $\left\langle\alpha_{i}, \alpha_{i}\right\rangle$ the length of any real root is equal to the length of some fundamental
root. The relative lengths of the fundamental roots may be obtained from list 17.1 using the formulae

$$
\begin{aligned}
& A_{i j}=2 \frac{\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} \\
& \frac{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}=\frac{A_{i j}}{A_{j i}} .
\end{aligned}
$$

Proposition 17.9 (a) If $A$ is an affine Cartan matrix of types $\tilde{A}_{l}, \tilde{D}_{l}, \tilde{E}_{6}, \tilde{E}_{7}, \tilde{E}_{8}$ all the fundamental roots have the same length.
(b) If $A$ has types $\tilde{B}_{l}, \tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}, \tilde{C}_{l}^{t}, \tilde{F}_{4}, \tilde{F}_{4}^{\mathrm{t}}$ there are fundamental roots of two different lengths. The ratio $\langle\beta, \beta\rangle /\langle\alpha, \alpha\rangle$ where $\alpha$ is short and $\beta$ is long is 2.
(c) If A has type $\tilde{G}_{2}$ or $\tilde{G}_{2}^{\mathrm{t}}$ there are fundamental roots of two different lengths with $\langle\beta, \beta\rangle /\langle\alpha, \alpha\rangle=3$.
(d) If $A$ has type $\tilde{A}_{1}^{\prime}$ there are fundamental roots of two different lengths with $\langle\beta, \beta\rangle /\langle\alpha, \alpha\rangle=4$.
(e) If A has type $\tilde{C}_{l}^{\prime}$ there are fundamental roots of three different lengths, say $\alpha, \beta, \gamma$, with $\langle\beta, \beta\rangle /\langle\alpha, \alpha\rangle=2$ and $\langle\gamma, \gamma\rangle /\langle\beta, \beta\rangle=2$.

Proof. This is clear from list 17.1.
We shall denote by $\Phi_{\mathrm{Re}, \mathrm{s}}$ the set of short real roots, by $\Phi_{\mathrm{Re}, 1}$ the set of long real roots and by $\Phi_{\mathrm{Re}, \mathrm{i}}$ the set of real roots of intermediate length. The latter set is non-empty only when $A$ has type $\tilde{C}_{l}^{\prime}$ for some $l$. If all real roots have the same length we use the convention $\Phi_{\mathrm{Re}}=\Phi_{\mathrm{Re}, \mathrm{s}}$.

We now aim to characterise the set $\Phi_{\mathrm{Re}, \mathrm{s}}$. We consider the possible values of $\langle\alpha, \alpha\rangle$ for $\alpha \in Q$. Let $\alpha=\sum_{i=0}^{l} k_{i} \alpha_{i}$. Then $\langle\alpha, \alpha\rangle=\sum_{i, j} k_{i} k_{j}\left\langle\alpha_{i}, \alpha_{j}\right\rangle$. Now $\left\langle\alpha_{i}, \alpha_{j}\right\rangle \in \mathbb{Q}$ for all $i, j$. Thus there exists $d \in \mathbb{Z}$ with $d>0$ such that $\left\langle\alpha_{i}, \alpha_{j}\right\rangle \in$ $\frac{1}{d} \mathbb{Z}$ for all $i, j$. Thus if $\langle\alpha, \alpha\rangle>0$ then $\langle\alpha, \alpha\rangle \geq \frac{1}{d}$. Hence there exists $m>0$ such that $m=\min \langle\alpha, \alpha\rangle$ for all $\alpha \in Q$ with $\langle\alpha, \alpha\rangle>0$.

Proposition 17.10 If $\alpha \in Q$ satisfies $\langle\alpha, \alpha\rangle=m$ then $\alpha \in Q^{+}$or $\alpha \in Q^{-}$.

Proof. Suppose if possible there exists $\alpha \in Q$ with $\langle\alpha, \alpha\rangle=m$ but $\alpha \notin Q^{+}$ and $\alpha \notin Q^{-}$. Then $\alpha=\beta-\gamma$ where $\beta, \gamma \in Q^{+}, \beta \neq 0, \gamma \neq 0$ and supp $\beta \cap$ supp $\gamma=\phi$. Hence

$$
\langle\alpha, \alpha\rangle=\langle\beta, \beta\rangle+\langle\gamma, \gamma\rangle-2\langle\beta, \gamma\rangle
$$

and $\langle\beta, \gamma\rangle \leq 0$ since supp $\beta \cap \operatorname{supp} \gamma=\phi$. Hence $\langle\alpha, \alpha\rangle \geq\langle\beta, \beta\rangle+\langle\gamma, \gamma\rangle$.

Now all proper connected principal minors of $A$ have finite type. Thus, considering the connected components of $\operatorname{supp} \beta$, we have $\beta=\beta_{1}+\cdots+\beta_{r}$ with supp $\beta_{i}$ connected for each $i,\left\langle\beta_{i}, \beta_{j}\right\rangle=0$ for $i \neq j$, and $\left\langle\beta_{i}, \beta_{i}\right\rangle>0$. Thus

$$
\langle\beta, \beta\rangle=\left\langle\beta_{1}, \beta_{1}\right\rangle+\cdots+\left\langle\beta_{r}, \beta_{r}\right\rangle>0 .
$$

Hence $\langle\beta, \beta\rangle \geq m$. Similarly $\langle\gamma, \gamma\rangle \geq m$. But then $\langle\alpha, \alpha\rangle \geq 2 m$, a contradiction. Hence $\alpha \in Q^{+}$or $\alpha \in Q^{-}$.

Proposition 17.11 Let $A$ be an indecomposable GCM of finite or affine type. Then the set $\Phi_{\mathrm{Re}, \mathrm{s}}$ of short real roots of $L(A)$ is given by

$$
\Phi_{\mathrm{re}, \mathrm{~s}}=\{\alpha \in Q ;\langle\alpha, \alpha\rangle=m\}
$$

Proof. Suppose $\alpha \in Q$ satisfies $\langle\alpha, \alpha\rangle=m$. We show $\alpha \in \Phi_{\mathrm{Re}}$. By Proposition $17.10 \alpha \in Q^{+}$or $\alpha \in Q^{-}$. We may suppose $\alpha \in Q^{+}$. Consider the set

$$
\{w(\alpha) ; w \in W\} \cap Q^{+}
$$

We choose an element $\beta=\sum k_{i} \alpha_{i}$ in this set with ht $\beta$ minimal. Then $\langle\beta, \beta\rangle=m$, so $\sum_{i} k_{i}\left\langle\alpha_{i}, \beta\right\rangle=m$. Since $k_{i} \geq 0$ and $m>0$ there exists $i$ with $\left\langle\alpha_{i}, \beta\right\rangle>0$. Thus $\beta\left(h_{i}\right)=2 \frac{\left\langle\alpha_{i}, \beta\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}>0$. Now $s_{i}(\beta)=\beta-\beta\left(h_{i}\right) \alpha_{i}$ so ht $s_{i}(\beta)<$ ht $\beta$. By minimality of ht $\beta$ we must have $s_{i}(\beta) \in Q^{-}$. But $\beta \in Q^{+}, s_{i}(\beta) \in Q^{-}$ imply $\beta=r \alpha_{i}$ for some $r \in \mathbb{Z}$ with $r>0$. Since

$$
\left\langle r \alpha_{i}, r \alpha_{i}\right\rangle=r^{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle \geq r^{2} m
$$

we have $r=1$. Thus $\beta=\alpha_{i}$ and $\left\langle\alpha_{i}, \alpha_{i}\right\rangle=m$. Hence $\beta \in \Phi_{\mathrm{Re}, \mathrm{s}}$ and so $\alpha \in \Phi_{\mathrm{Re}, \mathrm{s}}$ also.

Conversely if $\alpha \in \Phi_{\mathrm{Re}, \mathrm{s}}$ then $\alpha=w\left(\alpha_{i}\right)$ for some $w \in W$ and some $i$, and $\langle\alpha, \alpha\rangle=\left\langle\alpha_{i}, \alpha_{i}\right\rangle$. However, we have seen that the short fundamental roots have $\left\langle\alpha_{i}, \alpha_{i}\right\rangle=m$. Thus $\langle\alpha, \alpha\rangle=m$ also.

We aim next to characterise the set $\Phi_{\mathrm{Re}, 1}$ of long real roots. In order to do this we compare the roots of $L(A)$ and $L\left(A^{\mathrm{t}}\right)$. Here $A$ can be any GCM.

Proposition 17.12 If $\left(H, \Pi, \Pi^{v}\right)$ is a minimal realisation of the $G C M$ A then $\left(H^{*}, \Pi^{\mathrm{v}}, \Pi\right)$ is a minimal realisation of $A^{\mathrm{t}}$.

Proof. Let $A$ be an $n \times n$ matrix of rank $l$. Then $\operatorname{dim} H=2 n-l, \Pi^{\mathrm{v}}=$ $\left\{h_{1}, \ldots, h_{n}\right\}$ is a linearly independent subset of $H, \Pi=\left\{\alpha_{1}, \ldots, \alpha_{n}\right\}$ is a linearly independent subset of $H^{*}$, and $\alpha_{j}\left(h_{i}\right)=A_{i j}$.

We now replace $H$ by its dual space $H^{*}$. We still have $\operatorname{dim} H^{*}=2 n-l$, $\left(H^{*}\right)^{*}$ can be identified with $H$ by means of the formula

$$
h(\lambda)=\lambda(h) \quad \text { for } h \in H, \lambda \in H^{*}
$$

Since

$$
h_{j}\left(\alpha_{i}\right)=\alpha_{i}\left(h_{j}\right)=A_{j i}
$$

we see that $\left(H^{*}, \Pi^{\mathrm{v}}, \Pi\right)$ is a minimal realisation of $A^{\mathrm{t}}$.
Now suppose $A$ is symmetrisable. Then we have an isomorphism between $H$ and $H^{*}$ induced by our standard invariant form. Under this isomorphism $h_{i}$ corresponds to $\frac{2 \alpha_{i}}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}=d_{i} \alpha_{i}$. For each real root $\alpha \in \Phi_{\operatorname{Re}}$ we define the corresponding coroot $h_{\alpha} \in H$ to be the element of $H$ corresponding to $\frac{2 \alpha}{\langle\alpha, \alpha\rangle} \in$ $H^{*}$. The element $h_{\alpha}$ can also be described by using the Weyl group. Since the $W$-actions on $H$ and $H^{*}$ are compatible with the above isomorphism, if $\alpha=w\left(\alpha_{i}\right)$ then $h_{\alpha}=w\left(h_{i}\right)$. For $h_{i}$ corresponds to $\frac{2 \alpha_{i}}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}$ and $\langle\alpha, \alpha\rangle=\left\langle\alpha_{i}, \alpha_{i}\right\rangle$. Thus the coroots $h_{\alpha}$ for $\alpha \in \Phi_{\mathrm{Re}}$ for $L(A)$ may be interpreted as the real roots for $L\left(A^{\mathrm{t}}\right)$. Moreover we have

$$
\left\langle h_{\alpha}, h_{\alpha}\right\rangle=\left\langle\frac{2 \alpha}{\langle\alpha, \alpha\rangle}, \frac{2 \alpha}{\langle\alpha, \alpha\rangle}\right\rangle=\frac{4}{\langle\alpha, \alpha\rangle} .
$$

Hence $\alpha$ is a short root for $L(A)$ if and only if $h_{\alpha}$ is a long root for $L\left(A^{\mathrm{t}}\right)$. The fact that short roots give long coroots and long roots give short coroots is very useful. We shall apply this to characterise $\Phi_{\mathrm{Re}, 1}$ in the case when $A$ is of finite or affine type.

Proposition 17.13 Let A be an indecomposable GCM of finite or affine type. Then the set $\Phi_{\mathrm{Re}, 1}$ of long real roots of $L(A)$ is given by

$$
\Phi_{\mathrm{Re}, \mathrm{l}}=\left\{\alpha=\sum k_{i} \alpha_{i} \in Q ;\langle\alpha, \alpha\rangle=M, k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z} \quad \text { for all } i\right\}
$$

where $M=\max \left\{\langle\alpha, \alpha\rangle ; \alpha \in \Phi_{\mathrm{Re}}\right\}$.

Proof. We first show the long real roots satisfy the given conditions. Let $\alpha \in \Phi_{\mathrm{Re}, 1}$. Then $\langle\alpha, \alpha\rangle=M$. Let $\alpha=\sum k_{i} \alpha_{i}$. Then

$$
\frac{2 \alpha}{\langle\alpha, \alpha\rangle}=\sum k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle} \frac{2 \alpha_{i}}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}
$$

and so

$$
h_{\alpha}=\sum k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle} h_{i}
$$

This expresses a root for $L\left(A^{\mathrm{t}}\right)$ as a linear combination of fundamental roots, thus the coefficients $k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle}$ lie in $\mathbb{Z}$.

Conversely suppose $\alpha \in Q$ satisfies the given conditions. Then $h_{\alpha} \in \sum \mathbb{Z} h_{i}$ and $\left\langle h_{\alpha}, h_{\alpha}\right\rangle=4 / M$. Now $4 / M$ is the minimum possible value of $\langle\beta, \beta\rangle$ for all real roots $\beta$ of $L\left(A^{\mathrm{t}}\right)$. Thus by Proposition $17.11 h_{\alpha}$ is a short root of $L\left(A^{\mathrm{t}}\right)$. Hence $\alpha$ is a long root of $L(A)$.

We next wish to characterise the set $\Phi_{\mathrm{Re}, \mathrm{i}}$ of intermediate roots of $L(A)$ when $A$ has type $\tilde{C}_{l}^{\prime}$. We first need a lemma.

Lemma 17.14 (a) Suppose $A$ is an indecomposable GCM of finite or affine
 invariant under $W$.
(b) If $\alpha=\sum k_{i} \alpha_{i} \in Q$ satisfies $k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$ for all $i$ then $\alpha \in Q^{+}$or $\alpha \in Q^{-}$.

Proof. (a) Suppose $\alpha$ satisfies our condition. It is sufficient to show that $s_{j}(\alpha)$ satisfies it also. Now $s_{j}(\alpha)=\alpha-\alpha\left(h_{j}\right) \alpha_{j}$. Thus it is sufficient to show that

$$
\left(k_{j}-\alpha\left(h_{j}\right)\right) \frac{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}
$$

that is $\alpha\left(h_{j}\right) \frac{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$. Now we have

$$
\begin{aligned}
\alpha\left(h_{j}\right) \frac{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}{\langle\alpha, \alpha\rangle} & =\sum_{i} k_{i} \alpha_{i}\left(h_{j}\right) \frac{\left\langle\alpha_{j}, \alpha_{j}\right\rangle}{\langle\alpha, \alpha\rangle} \\
& =\sum_{i} k_{i} \alpha_{j}\left(h_{i}\right) \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle}=\sum_{i} A_{i j} k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}
\end{aligned}
$$

as required.
(b) Suppose the result is false. Then $\alpha=\beta-\gamma$ where $\beta, \gamma \in Q^{+}, \beta \neq 0, \gamma \neq 0$ and supp $\beta \cap \operatorname{supp} \gamma=\phi$. Then

$$
\langle\alpha, \alpha\rangle=\langle\beta, \beta\rangle+\langle\gamma, \gamma\rangle-2\langle\beta, \gamma\rangle \geq\langle\beta, \beta\rangle+\langle\gamma, \gamma\rangle .
$$

Now $\beta=\sum_{i \in \operatorname{supp} \beta} k_{i} \alpha_{i}$ so

$$
\langle\beta, \beta\rangle=\sum_{i} k_{i}^{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle+\sum_{i<j} 2 k_{i} k_{j}\left\langle\alpha_{i}, \alpha_{j}\right\rangle
$$

and

$$
\frac{\langle\beta, \beta\rangle}{\langle\alpha, \alpha\rangle}=\sum_{i} k_{i}\left(k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle}\right)+\sum_{i<j} A_{i j} k_{j}\left(k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle}\right) .
$$

Hence $\frac{\langle\beta, \beta\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$. Similarly we have $\frac{\langle\gamma, \gamma\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$.
Now all proper connected principal minors of $A$ have finite type. Thus we have $\beta=\beta_{1}+\cdots+\beta_{r}$ with supp $\beta_{i}$ connected for each $i,\left\langle\beta_{i}, \beta_{j}\right\rangle=0$ for $i \neq j$, and $\left\langle\beta_{i}, \beta_{i}\right\rangle>0$. Thus

$$
\langle\beta, \beta\rangle=\sum_{i}\left\langle\beta_{i}, \beta_{i}\right\rangle>0
$$

Similarly we can show $\langle\gamma, \gamma\rangle>0$. Thus $\langle\alpha, \alpha\rangle>0$ also. But now we have $\frac{\langle\beta, \beta\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$ so $\langle\beta, \beta\rangle \geq\langle\alpha, \alpha\rangle$, and $\frac{\langle\gamma, \gamma\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$ so $\langle\gamma, \gamma\rangle \geq\langle\alpha, \alpha\rangle$. Hence $\langle\alpha, \alpha\rangle \geq\langle\beta, \beta\rangle+\langle\gamma, \gamma\rangle \geq 2\langle\alpha, \alpha\rangle$, a contradiction.

We now suppose $A$ is a GCM of affine type $\tilde{C}_{l}^{\prime}$. The diagram of $A$ is


Let $m^{\prime}$ be defined by $\left\langle\alpha_{i}, \alpha_{i}\right\rangle=m^{\prime}$ for $i=1, \ldots, l-1$. Thus $m^{\prime}$ is the squared length of the intermediate roots.

Lemma 17.15 Suppose $A$ has type $\tilde{C}_{l}^{\prime}$. Suppose $\alpha=\sum_{i=0}^{l} k_{i} \alpha_{i} \in Q$ satisfies $\langle\alpha, \alpha\rangle=m^{\prime}$. Then $k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$ for all $i$.

Proof. The required condition is obvious for all $i \neq 0$ since $\left\langle\alpha_{i}, \alpha_{i}\right\rangle=m^{\prime}$ for $i=1, \ldots, l-1$ and $\left\langle\alpha_{l}, \alpha_{l}\right\rangle=2 m^{\prime}$. We must therefore show $k_{0} \frac{\left\langle\alpha_{0}, \alpha_{0}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$, that is that $k_{0}$ is even.

Now $\alpha=k_{0} \alpha_{0}+\sum_{i=1}^{l} k_{i} \alpha_{i}$, thus

$$
\begin{aligned}
\langle\alpha, \alpha\rangle & =k_{0}^{2}\left\langle\alpha_{0}, \alpha_{0}\right\rangle+2 k_{0} k_{1}\left\langle\alpha_{0}, \alpha_{1}\right\rangle+\left\langle\sum_{i=1}^{l} k_{i} \alpha_{i}, \sum_{i=1}^{l} k_{i} \alpha_{i}\right\rangle \\
& =k_{0}^{2}\left\langle\alpha_{0}, \alpha_{0}\right\rangle+k_{0} k_{1} A_{10}\left\langle\alpha_{1}, \alpha_{1}\right\rangle+\sum_{i=1}^{l} k_{i}^{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle+\sum_{\substack{i, j=1 \\
i<j}}^{l} k_{i} k_{j} A_{i j}\left\langle\alpha_{i}, \alpha_{i}\right\rangle .
\end{aligned}
$$

Thus $\langle\alpha, \alpha\rangle \in k_{0}^{2}\left\langle\alpha_{0}, \alpha_{0}\right\rangle+\mathbb{Z} m^{\prime}$. But $\langle\alpha, \alpha\rangle=m^{\prime}$ so $k_{0}^{2}\left\langle\alpha_{0}, \alpha_{0}\right\rangle \in \mathbb{Z} m^{\prime}$. Since $\left\langle\alpha_{0}, \alpha_{0}\right\rangle=\frac{1}{2} m^{\prime}$ we have $k_{0}^{2} / 2 \in \mathbb{Z}$ and so $k_{0}$ is even, as required.

We can now characterise $\Phi_{\mathrm{Re}, \mathrm{i}}$.
Proposition 17.16 Suppose $A$ is a GCM of type $\tilde{C}_{l}^{\prime}$. Then

$$
\Phi_{\mathrm{Re}, \mathrm{i}}=\left\{\alpha \in Q ;\langle\alpha, \alpha\rangle=m^{\prime}\right\} .
$$

Proof. Let $\alpha \in Q$ satisfy $\langle\alpha, \alpha\rangle=m^{\prime}$. By Lemma $17.15 \alpha=\sum_{i=0}^{l} k_{i} \alpha_{i}$ with $k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$ for each $i$. By Lemma 17.14(b) $\alpha \in Q^{+}$or $\alpha \in Q^{-}$. We may assume $\alpha \in Q^{+}$.

Consider the set

$$
\{w(\alpha) ; w \in W\} \cap Q^{+} .
$$

We choose an element $\beta=\sum_{i=0}^{l} k_{i}^{\prime} \alpha_{i}$ in this set with ht $\beta$ minimal. Then $\langle\beta, \beta\rangle=m^{\prime}$ and so $\sum_{i=0}^{l} k_{i}^{\prime}\left\langle\alpha_{i}, \beta\right\rangle=m^{\prime}$. Since $m^{\prime}>0$ and $k_{i}^{\prime} \geq 0$ there exists $i$ with $\left\langle\alpha_{i}, \beta\right\rangle>0$. Thus $\left.\beta\left(h_{i}\right)=2 \frac{\left\langle\alpha_{i}, \beta\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}\right\rangle 0$. Now $s_{i}(\beta)=\beta-\beta\left(h_{i}\right) \alpha_{i}$ so ht $s_{i}(\beta)<$ ht $\beta$. By the minimality of ht $\beta, s_{i}(\beta) \notin Q^{+}$. But $s_{i}(\beta) \in Q^{+}$or $Q^{-}$ by Lemma 17.14 (a) and (b). Thus $\beta \in Q^{+}$and $s_{i}(\beta) \in Q^{-}$. Hence $\beta=r \alpha_{i}$ for some $r \in \mathbb{Z}$ with $r>0$. Thus $\langle\beta, \beta\rangle=r^{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle=m^{\prime}$. However, $\left\langle\alpha_{i}, \alpha_{i}\right\rangle \geq$ $\frac{1}{2} m^{\prime}$ thus $r=1$. Thus $\beta=\alpha_{i} \in \Phi_{\mathrm{Re}, \mathrm{i}}$. It follows that $\alpha \in \Phi_{\mathrm{Re}, \mathrm{i}}$ also.
We are now able to obtain explicitly the set $\Phi_{\mathrm{Re}}$ of all real roots of each affine Kac-Moody algebra individually. We recall that $\Phi^{0}$ is the root system of the Lie algebra $L\left(A^{0}\right)$ of finite type obtained by removing vertex 0 from the diagram of $A$. We denote by $\Phi_{\mathrm{s}}^{0}$, $\Phi_{1}^{0}$ the set of short and long roots in $\Phi^{0}$. If all roots of $\Phi^{0}$ have the same length we write $\Phi_{\mathrm{s}}^{0}=\Phi^{0}$.

Theorem 17.17 The real roots of the affine Kac-Moody algebra $L(A)$ are as follows.
(a) If $A$ is one of the types $\tilde{A}_{l}, \tilde{B}_{l}, \tilde{C}_{l}, \tilde{D}_{l}, \tilde{E}_{6}, \tilde{E}_{7}, \tilde{E}_{8}, \tilde{F}_{4}, \tilde{G}_{2}$ then $\Phi_{\mathrm{Re}}=$ $\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\}$.
(b) If $A$ is one of the types $\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}$ then

$$
\begin{aligned}
& \Phi_{\mathrm{Re}, \mathrm{~s}}=\left\{\alpha+r \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, 1}=\left\{\alpha+2 r \delta ; \alpha \in \Phi_{1}^{0}, r \in \mathbb{Z}\right\} .
\end{aligned}
$$

(c) If A is of type $\tilde{G}_{2}^{t}$ then

$$
\begin{aligned}
& \Phi_{\mathrm{Re}, \mathrm{~s}}=\left\{\alpha+r \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, 1}=\left\{\alpha+3 r \delta ; \alpha \in \Phi_{1}^{0}, r \in \mathbb{Z}\right\} .
\end{aligned}
$$

(d) If $A$ is of type $\tilde{C}_{l}^{\prime}$ then

$$
\begin{aligned}
& \Phi_{\mathrm{Re}, \mathrm{~s}}=\left\{\frac{1}{2}(\alpha+(2 r-1) \delta) ; \alpha \in \Phi_{1}^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, \mathrm{i}}=\left\{\alpha+r \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, 1}=\left\{\alpha+2 r \delta ; \alpha \in \Phi_{1}^{0}, r \in \mathbb{Z}\right\} .
\end{aligned}
$$

(e) If $A$ is of type $\tilde{A}_{1}^{\prime}$ then

$$
\begin{aligned}
& \Phi_{\mathrm{Re}, \mathrm{~s}}=\left\{\frac{1}{2}(\alpha+(2 r-1) \delta) ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, \mathrm{l}}=\left\{\alpha+2 r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\}
\end{aligned}
$$

Proof. (i) Suppose first that $A$ is not of type $\tilde{C}_{l}^{\prime}$ or $\tilde{A}_{1}^{\prime}$. Then $\Phi_{\mathrm{s}}^{0} \subset \Phi_{\mathrm{Re}, \mathrm{s}}$. Let $\alpha \in \Phi_{\mathrm{s}}^{0}$. Then $\langle\alpha, \alpha\rangle=m$. Hence for $r \in \mathbb{Z}$ we have $\langle\alpha+r \delta, \alpha+r \delta\rangle=m$ since $\langle\alpha, \delta\rangle=0$ and $\langle\delta, \delta\rangle=0$. By Proposition 17.11 this implies $\alpha+r \delta \in \Phi_{\mathrm{Re}, \mathrm{s}}$.

Conversely suppose $\beta=\sum_{i=0}^{l} k_{i} \alpha_{i} \in \Phi_{\mathrm{Re}, \mathrm{s}}$. We have $a_{0}=1$, thus $\delta=\alpha_{0}+$ $\sum_{i=1}^{l} a_{i} \alpha_{i}$. Hence $\beta-k_{0} \delta=\sum_{i=1}^{l}\left(k_{i}-k_{0} a_{i}\right) \alpha_{i}$. Thus $\left\langle\beta-k_{0} \delta, \beta-k_{0} \delta\right\rangle=$ $\langle\beta, \beta\rangle=m$. Again by Proposition 17.11 we deduce $\beta-k_{0} \delta \in \Phi_{\mathrm{s}}^{0}$. Let $\alpha=\beta-k_{0} \delta$. Then $\beta=\alpha+k_{0} \delta$ for $\alpha \in \Phi_{\mathrm{s}}^{0}, k_{0} \in \mathbb{Z}$.

Thus the short roots in $\Phi$ have the required form. We now consider the long roots. We have $\Phi_{1}^{0} \subset \Phi_{\mathrm{Re}, 1}$.

Let $\alpha \in \Phi_{1}^{0}$. Then $\langle\alpha, \alpha\rangle=M$ and so $\langle\alpha+s \delta, \alpha+s \delta\rangle=M$ for all $s \in \mathbb{Z}$. Let $\alpha=\sum_{i=1}^{l} k_{i} \alpha_{i}$. By Proposition 17.13 we have $k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$ for $i=1, \ldots, l$. The same proposition shows that $\alpha+s \delta \in \Phi_{\mathrm{Re}, 1}$ if and only if $s a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$ for $i=0,1, \ldots, l$. Now $\left\langle\alpha_{i}, \alpha_{i}\right\rangle=\frac{2 c_{i}}{a_{i}}$, thus the condition is $\frac{2 c_{i}}{\langle\alpha, \alpha\rangle} s \in \mathbb{Z}$ for $i=0,1, \ldots, l$. We note that $\left\langle\alpha_{0}, \alpha_{0}\right\rangle=2$ since $a_{0}=1$.

First suppose that $\alpha_{0}$ is a long root, that is that we are in case (a). Then $\langle\alpha, \alpha\rangle=2$ and so $\frac{2 c_{i} s}{\langle\alpha, \alpha\rangle}=c_{i} s \in \mathbb{Z}$. Hence $\alpha+s \delta \in \Phi_{\mathrm{Re}, 1}$ for all $s \in \mathbb{Z}$.

Conversely suppose $\beta=\sum_{i=0}^{l} k_{i} \alpha_{i} \in \Phi_{\mathrm{Re}, 1}$. Then

$$
\beta-k_{0} \delta=\sum_{i=1}^{l}\left(k_{i}-k_{0} a_{i}\right) \alpha_{i}
$$

and we have $\left\langle\beta-k_{0} \delta, \beta-k_{0} \delta\right\rangle=\langle\beta, \beta\rangle=M$. Since $\beta \in \Phi_{\mathrm{Re}, 1}$ we have $k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\beta, \beta\rangle} \in \mathbb{Z}$ for $i=0,1, \ldots, l$. We have $k_{0} a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\beta, \beta\rangle} \in \mathbb{Z}$ also since $\left\langle\alpha_{i}, \alpha_{i}\right\rangle=\frac{2 c_{i}}{a_{i}}$ and $\langle\beta, \beta\rangle=2$. Hence $\beta-k_{0} \delta \in \Phi_{1}^{0}$ by Proposition 17.13. Thus $\beta=\alpha+k_{0} \delta$ for some $\alpha \in \Phi_{l}^{0}$ and $k_{0} \in \mathbb{Z}$.

Next suppose that $\alpha_{0}$ is a short root, i.e. that we are in case (b) or (c). Let $\frac{\langle\alpha, \alpha\rangle}{\left\langle\alpha_{0}, \alpha_{0}\right\rangle}=p$. Then $p=2$ in case (b) and $p=3$ in case (c). Thus

$$
\frac{2 c_{i}}{\langle\alpha, \alpha\rangle} s=c_{i} \frac{s}{p} \quad \text { since }\langle\alpha, \alpha\rangle=2 p
$$

Since $c_{0}=1$ this lies in $\mathbb{Z}$ for all $i=0,1, \ldots, l$ if and only if $s$ is divisible by $p$. Thus by Proposition $17.13 \alpha+p r \delta \in \Phi_{\mathrm{Re}, 1}$ for all $r \in \mathbb{Z}$.

Conversely suppose $\beta=\sum_{i=0}^{l} k_{i} \alpha_{i} \in \Phi_{\mathrm{Re}, 1}$. Then

$$
\beta-k_{0} \delta=\sum_{i=1}^{l}\left(k_{i}-k_{0} a_{i}\right) \alpha_{i} .
$$

We have $\left\langle\beta-k_{0} \delta, \beta-k_{0} \delta\right\rangle=\langle\beta, \beta\rangle=M$. Since $\beta \in \Phi_{\mathrm{Re}, 1}$ we have $k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\beta, \beta\rangle} \in \mathbb{Z}$ for $i=0,1, \ldots, l$. In particular $k_{0} \frac{\left\langle\alpha_{0}, \alpha_{0}\right\rangle}{\langle\beta, \beta\rangle}=\frac{k_{0}}{p} \in \mathbb{Z}$. We show $k_{0} a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\beta, \beta\rangle} \in \mathbb{Z}$ for $i=1, \ldots, l$. For $k_{0} a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\beta, \beta\rangle}=\frac{k_{0}}{p} c_{i} \in \mathbb{Z}$ since $\left\langle\alpha_{i}, \alpha_{i}\right\rangle=\frac{2 c_{i}}{a_{i}}$ and $\langle\beta, \beta\rangle=2 p$. Thus by Proposition $17.13 \beta-k_{0} \delta \in \Phi_{l}^{0}$. Let $\alpha=\beta-k_{0} \delta$. Then $\beta=\alpha+p r \delta$ for some $\alpha \in \Phi_{l}^{0}, r \in \mathbb{Z}$.

We have thus proved the required result in cases (a), (b) and (c).
(ii) We now suppose that $A$ has type $\tilde{C}_{l}^{\prime}$. Then we have $\Phi_{s}^{0} \subset \Phi_{\mathrm{Re}, \mathrm{i}}$ and $\Phi_{l}^{0} \subset$ $\Phi_{\mathrm{Re}, 1}$. First suppose $\alpha \in \Phi_{s}^{0}$. Then $\langle\alpha, \alpha\rangle=m^{\prime}$ and so $\langle\alpha+r \delta, \alpha+r \delta\rangle=m^{\prime}$. By Proposition $17.16 \alpha+r \delta \in \Phi_{\mathrm{Re}, \mathrm{i}}$ for all $\alpha \in \Phi_{s}^{0}, r \in \mathbb{Z}$.

Conversely suppose $\beta=\sum_{i=0}^{l} k_{i} \alpha_{i} \in \Phi_{\mathrm{Re}, \mathrm{i}}$. Then $k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\beta, \beta\rangle} \in \mathbb{Z}$ for $i=0,1, \ldots, l$ by Lemma 17.15, in particular $k_{0} \frac{\left\langle\alpha_{0}, \alpha_{0}\right\rangle}{\langle\beta, \beta\rangle}=\frac{k_{0}}{2} \in \mathbb{Z}$. Now

$$
\beta-\frac{k_{0}}{2} \delta=\sum_{i=1}^{l}\left(k_{i}-\frac{k_{0}}{2} a_{i}\right) \alpha_{i} .
$$

We have $\left\langle\beta-\frac{k_{0}}{2} \delta, \beta-\frac{k_{0}}{2} \delta\right\rangle=m^{\prime}$ and so by Proposition $17.11 \beta-\frac{k_{0}}{2} \delta \in$ $\Phi_{s}^{0}$. Let $\alpha=\beta-\frac{k_{0}}{2} \delta$. Then $\beta=\alpha+\frac{k_{0}}{2} \delta=\alpha+r \delta$ for some $\alpha \in \Phi_{s}^{0}, r \in \mathbb{Z}$.

We now turn from the intermediate roots to the long roots. Suppose $\alpha \in \Phi_{l}^{0}$. Then $\langle\alpha, \alpha\rangle=M$ and $\langle\alpha+s \delta, \alpha+s \delta\rangle=M$ for $s \in \mathbb{Z}$. Let $\alpha=$ $\sum_{i=1}^{l} k_{i} \alpha_{i}$. Then $k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$ for $i=1, \ldots, l$. Now

$$
\alpha+s \delta=\sum_{i=1}^{l} k_{i} \alpha_{i}+\sum_{i=0}^{l} s a_{i} \alpha_{i}
$$

We wish to know for which $s \in \mathbb{Z}$ we have

$$
s a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z} \quad \text { for all } i=0,1, \ldots, l
$$

Now $\left\langle\alpha_{0}, \alpha_{0}\right\rangle=\frac{2 c_{0}}{a_{0}}=1$, thus $\langle\alpha, \alpha\rangle=4$. Hence $s a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle}=\frac{c_{i} S}{2}$. Since $c_{0}=1$ this lies in $\mathbb{Z}$ for all $i=0,1, \ldots, l$ if and only if $s$ is even. Thus by Proposition $17.13 \alpha+2 r \delta \in \Phi_{\mathrm{Re}, 1}$ for all $\alpha \in \Phi_{l}^{0}, r \in \mathbb{Z}$.

Conversely suppose $\beta=\sum_{i=0}^{l} k_{i} \alpha_{i} \in \Phi_{\mathrm{Re}, 1}$. Then $k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\beta, \beta\rangle} \in \mathbb{Z}$ for $i=$ $0,1, \ldots, l$, in particular $k_{0} \frac{\left\langle\alpha_{0}, \alpha_{0}\right\rangle}{\langle\beta, \beta\rangle}=\frac{k_{0}}{4} \in \mathbb{Z}$. Now

$$
\beta-\frac{k_{0}}{2} \delta=\sum_{i=1}^{l}\left(k_{i}-\frac{k_{0}}{2} a_{i}\right) \alpha_{i}
$$

satisfies $\left\langle\beta-\frac{k_{0}}{2} \delta, \beta-\frac{k_{0}}{2} \delta\right\rangle=M$. Also $\frac{k_{0}}{2} a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\beta, \beta\rangle}=\frac{k_{0}}{4} c_{i} \in \mathbb{Z}$. Thus by Proposition 17.13 we have $\beta-\frac{k_{0}}{2} \delta \in \Phi_{l}^{0}$. Let $\alpha=\beta-\frac{k_{0}}{2} \delta$. Then $\beta=$ $\alpha+2 r \delta$ for some $\alpha \in \Phi_{l}^{0}, r \in \mathbb{Z}$.

We now consider the short roots of $\Phi$. There is no root of $\Phi^{0}$ of the same length as the short roots of $\Phi_{\mathrm{Re}}$. The squared length of the short roots of $\Phi_{\mathrm{Re}}$ is one half that of the long roots of $\Phi^{0}$. So suppose $\alpha \in \Phi_{l}^{0}$. We consider elements of form $\frac{1}{2}(\alpha+s \delta)$ where $s \in \mathbb{Z}$. We consider which of these elements lie in $Q$. Since the long roots of $\Phi^{0}$ have form

$$
\pm\left\{\alpha_{l}, 2 \alpha_{l-1}+\alpha_{l}, \ldots, 2 \alpha_{1}+\cdots+2 \alpha_{l-1}+\alpha_{l}\right\}
$$

and $\delta=2 \alpha_{0}+2 \alpha_{1}+\cdots+2 \alpha_{l-1}+\alpha_{l}$ we see that $\frac{1}{2}(\alpha+s \delta) \in Q$ if and only if $s$ is odd. Thus we consider elements of $Q$ of form $\frac{1}{2}(\alpha+(2 r-1) \delta)$ with $r \in \mathbb{Z}$. We have

$$
\left\langle\frac{1}{2}(\alpha+(2 r-1) \delta), \frac{1}{2}(\alpha+(2 r-1) \delta)\right\rangle=\frac{1}{4}\langle\alpha, \alpha\rangle=m
$$

By Proposition 17.11 this implies that $\frac{1}{2}(\alpha+(2 r-1) \delta) \in \Phi_{\mathrm{Re}, \mathrm{s}}$.
Conversely suppose $\beta=\sum_{i=0}^{l} k_{i} \alpha_{i} \in \Phi_{\mathrm{Re}, \mathrm{s}}$. Then $k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\beta, \beta\rangle} \in \mathbb{Z}$ for $i=$ $0,1, \ldots, l$. Then $2 \beta-k_{0} \delta=\sum_{i=1}^{l}\left(2 k_{i}-k_{0} a_{i}\right) \alpha_{i}$. We have $\left\langle 2 \beta-k_{0} \delta\right.$, $\left.2 \beta-k_{0} \delta\right\rangle=4\langle\beta, \beta\rangle=4$.

This is the squared length of the elements of $\Phi_{l}^{0}$. We also have

$$
\frac{k_{0} a_{i}\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle 2 \beta, 2 \beta\rangle}=\frac{k_{0}}{2} c_{i} \in \mathbb{Z} \quad \text { for } i=1, \ldots, l
$$

since $k_{0} \in \mathbb{Z}$ and $c_{i}=2$ for $i=1, \ldots, l$. By Proposition 17.13 we have $2 \beta-k_{0} \delta \in \Phi_{l}^{0}$. Let $\alpha=2 \beta-k_{0} \delta$. Then $\beta=\frac{1}{2}\left(\alpha+k_{0} \delta\right)$. Since $\beta \in Q$, $k_{0}$ is odd. Thus

$$
\beta=\frac{1}{2}(\alpha+(2 r-1) \delta) \quad \text { for some } \alpha \in \Phi_{l}^{0}, r \in \mathbb{Z}
$$

(iii) Finally we suppose that $A$ has type $\tilde{A}_{1}^{\prime}$. The diagram of $A$ is

with $a_{0}=2, a_{1}=1, c_{0}=1, c_{1}=2$. We also have

$$
\left\langle\alpha_{0}, \alpha_{0}\right\rangle=1,\left\langle\alpha_{1}, \alpha_{1}\right\rangle=4
$$

Now $\Phi^{0} \subset \Phi_{\text {Re, }, \text {. }}$ Let $\alpha \in \Phi^{0}$. Then $\langle\alpha+s \delta, \alpha+s \delta\rangle=M$. We can write $\alpha=k_{1} \alpha_{1}$. Then $\alpha+s \delta=2 s \alpha_{0}+\left(k_{1}+s\right) \alpha_{1}$. We have $k_{1} \frac{\left\langle\alpha_{1}, \alpha_{1}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$ and we consider which $s \in \mathbb{Z}$ have the property that $s a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle} \in \mathbb{Z}$ for $i=0,1$. Since $s a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle}=\frac{s c_{i}}{2}$ and $c_{0}=1$ this lies in $\mathbb{Z}$ for $i=0,1$ if and only if $s$ is even. By Proposition 17.13 we deduce that $\alpha+2 r \delta \in \Phi_{\mathrm{Re}, 1}$ for all $\alpha \in \Phi^{0}, r \in \mathbb{Z}$.

Conversely suppose $\beta=k_{0} \alpha_{0}+k_{1} \alpha_{1}$ lies in $\Phi_{\mathrm{Re}, 1}$. Then $\beta-\frac{k_{0}}{2} \delta=$ $\left(k_{1}-\frac{k_{0} a_{1}}{2}\right) \alpha_{1}$. We have $\left\langle\beta-\frac{k_{0}}{2} \delta, \beta-\frac{k_{0}}{2} \delta\right\rangle=M$. Also $k_{0} \frac{\left\langle\alpha_{0}, \alpha_{0}\right\rangle}{\langle\beta, \beta\rangle}=\frac{k_{0}}{4} \in \mathbb{Z}$. By Proposition 17.11 we have $\beta-\frac{k_{0}}{2} \delta \in \Phi^{0}$. Let $\alpha=\beta-\frac{k_{0}}{2} \delta$. Then $\beta=$ $\alpha+\frac{k_{0}}{2} \delta=\alpha+2 r \delta$ for some $\alpha \in \Phi^{0}, r \in \mathbb{Z}$.
We now consider the short roots. Suppose $\alpha \in \Phi^{0}$ and consider the element $\frac{1}{2}(\alpha+s \delta)$ for $s \in \mathbb{Z}$. Since $\alpha= \pm \alpha_{1}$ and $\delta=2 \alpha_{0}+\alpha_{1}$ this element lies in $Q$ if and only if $s$ is odd. We have

$$
\left\langle\frac{1}{2}(\alpha+(2 r-1) \delta), \frac{1}{2}(\alpha+(2 r-1) \delta)\right\rangle=\frac{1}{4}\langle\alpha, \alpha\rangle=1 .
$$

This is the squared length of the short roots of $\Phi_{\mathrm{Re}, s}$. By Proposition 17.11 $\frac{1}{2}(\alpha+(2 r-1) \delta) \in \Phi_{\text {Re,s }}$ for all $r \in \mathbb{Z}$.
Conversely suppose $\beta=k_{0} \alpha_{0}+k_{1} \alpha_{1} \in \Phi_{\mathrm{Re}, s}$. Then

$$
\beta-\frac{k_{0}}{2} \delta=\left(k_{1}-\frac{k_{0} a_{1}}{2}\right) \alpha_{1} .
$$

We have $\left\langle 2 \beta-k_{0} \delta, 2 \beta-k_{0} \delta\right\rangle=4\langle\beta, \beta\rangle=4$. This is the squared length of the roots in $\Phi^{0}$. So by Proposition 17.11 we have $2 \beta-k_{0} \delta \in \Phi^{0}$. Let $\alpha=2 \beta-k_{0} \delta$. Then $\beta=\frac{1}{2}\left(\alpha+k_{0} \delta\right)$. Since $\beta \in Q k_{0}$ must be odd. Hence $\beta=\frac{1}{2}(\alpha+(2 r-1) \delta)$ for some $\alpha \in \Phi^{0}, r \in \mathbb{Z}$. This completes the proof.

### 17.3 The Weyl group of an affine Kac-Moody algebra

Let $A$ be an affine Cartan matrix and $W$ the Weyl group of $L(A)$. Then $W=\left\langle s_{0}, s_{1}, \ldots, s_{l}\right\rangle$. The subgroup $W^{0}=\left\langle s_{1}, \ldots, s_{l}\right\rangle$ is the Weyl group of the finite dimensional simple Lie algebra $L\left(A^{0}\right)$. In order to investigate the structure of $W$ we introduce the element $\theta=\delta-a_{0} \alpha_{0}=\sum_{i=1}^{l} a_{i} \alpha_{i}$. This element $\theta$ lies in $Q^{0}=Q\left(A^{0}\right)$.

Proposition 17.18 (i) If the affine Cartan matrix $A$ is not of type $\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}}$ then $\theta$ is the highest root of $\Phi^{0}$.
(ii) If A is of type $\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}}$ then $\theta$ is the highest short root of $\Phi^{0}$.

Proof. We first show that $\theta \in \Phi^{0}$. We have $\langle\theta, \theta\rangle=\left\langle\delta-a_{0} \alpha_{0}, \delta-a_{0} \alpha_{0}\right\rangle=$ $a_{0}^{2}\left\langle\alpha_{0}, \alpha_{0}\right\rangle=2 a_{0}$. First suppose $a_{0}=1$ and $\alpha_{0}$ is a long root. Then $\langle\theta, \theta\rangle=$ $\left\langle\alpha_{0}, \alpha_{0}\right\rangle=2$. Also

$$
a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\theta, \theta\rangle}=c_{i} \in \mathbb{Z} .
$$

Thus $\theta \in \Phi_{l}^{0}$ by Proposition 17.13
Next suppose $a_{0}=2$. Then $\langle\theta, \theta\rangle=4\left\langle\alpha_{0}, \alpha_{0}\right\rangle=4$. Thus $\theta$ has the same squared length as a long root. Also $a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\theta, \theta\rangle}=\frac{c_{i}}{2}$. This lies in $\mathbb{Z}$ for $i=1, \ldots, l$ since $c_{i}=2$ for such values of $i$. Hence $\theta \in \Phi_{l}^{0}$ by Proposition 17.13.

Finally suppose $a_{0}=1$ and $\alpha_{0}$ is a short root. This occurs for the cases in (ii). Then $\langle\theta, \theta\rangle=\left\langle\alpha_{0}, \alpha_{0}\right\rangle$. Hence $\theta \in \Phi_{s}^{0}$ by Proposition 17.11.

Thus we have shown $\theta \in \Phi^{0}$ in all cases. We also have

$$
\begin{aligned}
\left\langle\theta, \alpha_{i}\right\rangle & =\left\langle\delta-a_{0} \alpha_{0}, \alpha_{i}\right\rangle=-a_{0}\left\langle\alpha_{0}, \alpha_{i}\right\rangle \\
& =-\frac{a_{0} A_{0 i}\left\langle\alpha_{0}, \alpha_{0}\right\rangle}{2}=-c_{0} A_{0 i}=-A_{0 i} .
\end{aligned}
$$

Thus $\left\langle\theta, \alpha_{i}\right\rangle \geq 0$ for $i=1, \ldots, l$. Hence $\theta \in \bar{C}_{0}$, the closure of the fundamental chamber for $\Phi^{0}$. This implies that $\theta$ is the highest root of $\Phi^{0}$ in the cases in (i) and the highest short root of $\Phi^{0}$ in the cases in (ii), by Proposition 12.9.

Now let $s_{\theta}$ be the reflection corresponding to the root $\theta$. Then $s_{\theta}: H^{0} \rightarrow H^{0}$ is given by $s_{\theta}(h)=h-\theta(h) h_{\theta}$.

Lemma 17.19 The coroot $h_{\theta}$ is given by $h_{\theta}=\frac{1}{a_{0}}\left(c-h_{0}\right)$.
Proof. Since $\theta=\sum_{i=1}^{l} a_{i} \alpha_{i}$ we have $\frac{2 \theta}{\langle\theta, \theta\rangle}=\sum_{i=1}^{l} a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\langle\theta\rangle} \frac{2 \alpha_{i}}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}$, hence

$$
\begin{aligned}
h_{\theta} & =\sum_{i=1}^{l} a_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\theta, \theta\rangle} h_{i}=\sum_{i=1}^{l} \frac{2 c_{i} h_{i}}{2 a_{0}} \\
& =\frac{1}{a_{0}} \sum_{i=1}^{l} c_{i} h_{i}=\frac{1}{a_{0}}\left(c-h_{0}\right) .
\end{aligned}
$$

Now the affine Weyl group $W$ is generated by $W^{0}$ and $s_{0}$, so is also generated by $W^{0}$ and $s_{0} s_{\theta}$. We consider the action of $s_{0} s_{\theta}$ on $H$.

Proposition 17.20 Let $h \in H$. Then

$$
s_{0} s_{\theta}(h)=h+\delta(h) h_{\theta}-\left(\left\langle h_{\theta}, h\right\rangle+\frac{1}{2}\left\langle h_{\theta}, h_{\theta}\right\rangle \delta(h)\right) c .
$$

Proof.

$$
\begin{aligned}
s_{0} s_{\theta}(h) & =s_{0}\left(h-\theta(h) h_{\theta}\right)=h-\alpha_{0}(h) h_{0}-\theta(h)\left(h_{\theta}-\alpha_{0}\left(h_{\theta}\right) h_{0}\right) \\
& =h-\alpha_{0}(h)\left(c-a_{0} h_{\theta}\right)-\theta(h) h_{\theta}+\theta(h) \alpha_{0}\left(h_{\theta}\right)\left(c-a_{0} h_{\theta}\right) \\
& =h+\left(a_{0} \alpha_{0}(h)-\theta(h)-a_{0} \theta(h) \alpha_{0}\left(h_{\theta}\right)\right) h_{\theta}+\left(\theta(h) \alpha_{0}\left(h_{\theta}\right)-\alpha_{0}(h)\right) c .
\end{aligned}
$$

Now $\alpha_{0}\left(h_{\theta}\right)=\alpha_{0}\left(\frac{1}{a_{0}}\left(c-h_{0}\right)\right)=-\frac{2}{a_{0}}$. Thus

$$
\begin{aligned}
s_{0} s_{\theta}(h) & =h+\left(a_{0} \alpha_{0}(h)+\theta(h)\right) h_{\theta}-\left(\frac{2}{a_{0}} \theta(h)+\alpha_{0}(h)\right) c \\
& =h+\delta(h) h_{\theta}-\frac{1}{a_{0}}(\theta(h)+\delta(h)) c \\
& =h+\delta(h) h_{\theta}-\left(\left\langle h_{\theta}, h\right\rangle+\frac{1}{2}\left\langle h_{\theta}, h_{\theta}\right\rangle \delta(h)\right) c
\end{aligned}
$$

since $\left\langle h_{\theta}, h\right\rangle=\frac{2 \theta(h)}{\langle\theta, \theta\rangle}=\frac{1}{a_{0}} \theta(h)$ and $\left\langle h_{\theta}, h_{\theta}\right\rangle=\left\langle\frac{2 \theta}{\langle\theta, \theta\rangle}, \frac{2 \theta}{\langle\theta, \theta\rangle}\right\rangle=\frac{4}{\langle\theta, \theta\rangle}=\frac{2}{a_{0}}$. We define $t_{h_{\theta}}: H \rightarrow H$ by

$$
t_{h_{\theta}}(h)=h+\delta(h) h_{\theta}-\left(\left\langle h_{\theta}, h\right\rangle+\frac{1}{2}\left\langle h_{\theta}, h_{\theta}\right\rangle \delta(h)\right) c .
$$

Thus we have $s_{0} s_{\theta}=t_{h_{\theta}}$. Hence $W$ is generated by $W^{0}$ and $t_{h_{\theta}}$.
More generally, for any $x \in H^{0}$ we define $t_{x}: H \rightarrow H$ by

$$
t_{x}(h)=h+\delta(h) x-\left(\langle x, h\rangle+\frac{1}{2}\langle x, x\rangle \delta(h)\right) c .
$$

Proposition 17.21 (i) $t_{x} t_{y}=t_{x+y}$ for all $x, y \in H^{0}$.
(ii) $w t_{x} w^{-1}=t_{w(x)}$ for all $w \in W^{0}, x \in H^{0}$.

Proof. The linear map $t_{x}: H \rightarrow H$ is uniquely determined by the properties

$$
\begin{aligned}
& t_{x}(h)=h-\langle x, h\rangle c \quad \text { when } \delta(h)=0 \\
& t_{x}(d)=d+a_{0} x-\frac{1}{2} a_{0}\langle x, x\rangle c
\end{aligned}
$$

since $\delta\left(h_{i}\right)=0$ and $\delta(d)=a_{0}$. If $\delta(h)=0$ then

$$
\begin{aligned}
t_{x} t_{y}(h) & =t_{x}(h-\langle y, h\rangle c)=h-\langle x, h\rangle c-\langle y, h\rangle(c-\langle x, c\rangle c) \\
& =h-\langle x+y, h\rangle c \quad \text { since }\langle x, c\rangle=0 \\
& =t_{x+y}(h)
\end{aligned}
$$

Also

$$
\begin{aligned}
t_{x} t_{y}(d) & =t_{x}\left(d+a_{0} y-\frac{1}{2} a_{0}\langle y, y\rangle c\right) \\
& =d+a_{0} x-\frac{1}{2} a_{0}\langle x, x\rangle c+a_{0}(y-\langle x, y\rangle c)-\frac{1}{2} a_{0}\langle y, y\rangle(c-\langle x, c\rangle c) \\
& =d+a_{0}(x+y)-a_{0}\left(\frac{1}{2}\langle x, x\rangle+\langle x, y\rangle+\frac{1}{2}\langle y, y\rangle\right) c \\
& =d+a_{0}(x+y)-a_{0} \cdot \frac{1}{2}\langle x+y, x+y\rangle c \\
& =t_{x+y}(d)
\end{aligned}
$$

Thus $t_{x} t_{y}=t_{x+y}$ for all $x, y \in H^{0}$.
Now let $w \in W^{0}$, and $h \in H$ satisfy $\delta(h)=0$. Then

$$
w t_{x} w^{-1}(h)=w\left(w^{-1}(h)-\left\langle x, w^{-1}(h)\right\rangle c\right)
$$

since $\delta\left(w^{-1}(h)\right)=(w \delta)(h)=\delta(h)=0$. Thus

$$
w t_{x} w^{-1}(h)=h-\langle w(x), h\rangle c=t_{w(x)}(h)
$$

since $w(c)=c$. Also $w(d)=d$ for all $w \in W^{0}$ and so

$$
\begin{aligned}
w t_{x} w^{-1}(d) & =w t_{x}(d)=w\left(d+a_{0} x-\frac{1}{2}\langle x, x\rangle a_{0} c\right) \\
& =d+a_{0} w(x)-\frac{1}{2}\langle w(x), w(x)\rangle a_{0} c \\
& =t_{w(x)}(d)
\end{aligned}
$$

Hence $w t_{x} w^{-1}=t_{w(x)}$.
Let $M$ be the additive subgroup (i.e. lattice) of $H^{0}$ generated by the elements $w\left(h_{\theta}\right)$ for all $w \in W^{0}$. Let $t(M)=\left\{t_{m} ; m \in M\right\}$.

Proposition 17.22 $W=t(M) W^{0}$ where $t(M)$ is normal in $W$ and $t(M) \cap W^{0}=$ 1. Thus $W$ is a semidirect product of $t(M)$ and $W^{0}$.

Proof. We know that $t_{h_{\theta}} \in W$, hence $w t_{h_{\theta}} w^{-1}=t_{w\left(h_{\theta}\right)} \in W$ for all $w \in W^{0}$. Thus $t(M)$ is a subgroup of $W$. Since $W$ is generated by $W^{0}$ and $t_{h_{\theta}}, W$ is generated by $t(M)$ and $W^{0}$. But $W^{0}$ lies in the normaliser of $t(M)$ by Proposition 17.21 (ii). Thus $W=t(M) W^{0}$. Finally $t(M) \cap W^{0}=1$ since $t(M)$ is a free abelian group whereas $W^{0}$ is finite.

The lattice $M \subset H_{\mathbb{R}}^{0}$ will be important in understanding the affine Weyl group $W$. We shall now identify it in each case.

Proposition 17.23 (i) If $A$ is an affine Cartan matrix not of types $\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}$, $\tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}}$ then $M=\sum_{i=1}^{l} \mathbb{Z} h_{i}$.
(ii) If A has type $\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}}$ then

$$
M=\sum_{\alpha_{i} \text { short }} \mathbb{Z} h_{i}+\sum_{\alpha_{i} \text { long }} p \mathbb{Z} h_{i}
$$

where $p$ is the ratio of the squared lengths of the long and short roots ( $p=3$ for $\tilde{G}_{2}^{t}$ and $p=2$ in the other cases).

Proof. By Proposition $17.18 \theta$ is a long root in the cases in (i) and a short root in the cases in (ii). Thus $h_{\theta}$ is a short coroot in (i) and a long coroot in (ii). Thus $M$ is generated by all short coroots in (i) and by all long coroots in (ii). Now it follows from Proposition 8.18 that the set of all short coroots generates the coroot lattice $\sum_{i=1}^{l} \mathbb{Z} h_{i}$. But the set of all long coroots generates the sublattice with basis $h_{i}$ for $h_{i}$ long (i.e. $\alpha_{i}$ short) and $p h_{i}$ for $h_{i}$ short (i.e. $\alpha_{i}$ long). The result follows.

We have been considering an action of the affine Weyl group $W$ by linear transformations of the vector space $H$ of dimension $l+2$. However, we now show that there is a simpler action of $W$ by affine transformations on the real vector space $H_{\mathbb{R}}^{0}$ of dimension $l$. We recall that the group of affine transformations of a vector space is generated by the group of non-singular linear transformations and the group of translations.

We first define $H_{\mathbb{R}, 1}=\left\{h \in H_{\mathbb{R}} ; \delta(h)=1\right\}$. The space $H_{\mathbb{R}, 1}$, although not a subspace of $H_{\mathbb{R}}$, is invariant under $W$. For

$$
\delta(w(h))=\left(w^{-1} \delta\right)(h)=\delta(h)
$$

since $w(\delta)=\delta$. Now we have a decomposition

$$
H_{\mathbb{R}}=H_{\mathbb{R}}^{0} \oplus(\mathbb{R} c+\mathbb{R} d)
$$

into subspaces of dimension $l$ and 2 which are mutually orthogonal. For $\left\langle h_{i}, c\right\rangle=0$ and $\left\langle h_{i}, d\right\rangle=0$ for $i=1, \ldots, l$. Since $\delta\left(h_{i}\right)=0$ for $i=$ $1, \ldots, l, \delta(c)=0$ and $\delta(d)=a_{0}$ the elements of $H_{\mathbb{R}}$ which lie in $H_{\mathbb{R}, 1}$ are those of form

$$
\sum_{i=1}^{l} \lambda_{i} h_{i}+\lambda c+\frac{1}{a_{0}} d \quad \lambda_{i} \in \mathbb{R}, \lambda \in \mathbb{R} .
$$

Now $h \in H_{\mathbb{R}, 1}$ implies $h+\mu c \in H_{\mathbb{R}, 1}$ for $\mu \in \mathbb{R}$. Since $w(c)=c$ for all $w \in$ $W, W$ acts on the quotient space $H_{\mathbb{R}, 1} / \mathbb{R} c$. Also we have a bijective map

$$
H_{\mathbb{R}, 1} / \mathbb{R} c \rightarrow H_{\mathbb{R}}^{0}
$$

given by

$$
\mathbb{R} c+\sum_{i=1}^{l} \lambda_{i} h_{i}+\frac{1}{a_{0}} d \rightarrow \sum_{i=1}^{l} \lambda_{i} h_{i}
$$

and this bijection may be used to define an action of $W$ on $H_{\mathbb{R}}^{0}$.
Proposition 17.24 The action of $W=t(M) W^{0}$ on $H_{\mathbb{R}}^{0}$ is as follows. The $W^{0}$-action on $H_{\mathbb{R}}^{0}$ is that previously considered. For $m \in M, h \in H_{\mathbb{R}}^{0}$ we have $t_{m}(h)=h+m$. Thus $t_{m}$ acts on $H_{\mathbb{R}}^{0}$ as translation by $m$. Hence $W$ acts on $H_{\mathbb{R}}^{0}$ as a group of affine transformations.

Proof. If $w \in W^{0}$ then $w(c)=c$ and $w(d)=d$. This implies that the $w$-action on $H_{\mathbb{R}}^{0}$ defined above is the usual $w$-action. If $m \in M, h \in H_{\mathbb{R}, 1}$ then $t_{m}(h)=$ $h+m+\mu c$ for some $\mu \in \mathbb{R}$. This induces an action of $t_{m}$ on $H_{\mathbb{R}}^{0}$ given by $t_{m}(h)=h+m$. Thus $t_{m}$ acts on $H_{\mathbb{R}}^{0}$ as translation by $m$.

Corollary 17.25 The action of $W$ on $H_{\mathbb{R}}^{0}$ is faithful.
Proof. Suppose $t_{m} w, w \in W^{0}$, acts trivially on $H_{\mathbb{R}}^{0}$. Then $t_{m} w(0)=0$. This implies $m=0$, that is $t_{m}=1$. Hence $w \in W^{0}$ acts trivially on $H_{\mathbb{R}}^{0}$. Since $W^{0}$ acts faithfully on $H_{\mathbb{R}}^{0}$ this implies $w=1$.

Corollary $\mathbf{1 7 . 2 6} s_{0}$ acts on $H_{\mathbb{R}}^{0}$ as the reflection in the affine hyperplane

$$
L_{\theta, 1}=\left\{h \in H_{\mathbb{R}}^{0} ; \quad \theta(h)=1\right\}
$$

Proof. For $h \in H_{\mathbb{R}}^{0}$ we have

$$
s_{0}(h)=t_{h_{\theta}} s_{\theta}(h)=h-\theta(h) h_{\theta}+h_{\theta}=h+(1-\theta(h)) h_{\theta} .
$$

This is the reflection in $L_{\theta, 1}$.
For each $\alpha \in \Phi^{0}$ and $k \in \mathbb{Z}$ let $L_{\alpha, k}$ be the affine hyperplane given by

$$
L_{\alpha, k}=\left\{h \in H_{\mathbb{R}}^{0} ; \alpha(h)=k\right\} .
$$

Thus the generators $s_{0}, s_{1}, \ldots, s_{l}$ of the affine Weyl group $W$ act on $H_{\mathbb{R}}^{0}$ as the reflections in the hyperplanes $L_{\theta, 1}, L_{\alpha_{1}, 0}, \ldots, L_{\alpha_{l}, 0}$ respectively.

We now introduce a collection of affine hyperplanes whose corresponding affine reflections will lie in $W$.

Let $\mathfrak{Z}=\left\{L_{\alpha, k} ; \alpha \in \Phi^{0}, k \in \mathbb{Z}, p\right.$ divides $k$ if $\alpha$ is a long root and $A$ is one of $\left.\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}}\right\}$. Here as usual $p=2$ in the first three cases and $p=3$ for $\tilde{G}_{2}^{\mathrm{t}}$.

Let $s_{\alpha, k}$ be the reflection in $L_{\alpha, k}$. Then $s_{\alpha, k}(h)=h+(k-\alpha(h)) h_{\alpha}$. For

$$
\frac{h+\left(h+(k-\alpha(h)) h_{\alpha}\right)}{2} \in L_{\alpha, k}
$$

and $h+(k-\alpha(h)) h_{\alpha}$ differs from $h$ by a multiple of $h_{\alpha}$. Thus $s_{\alpha, k}=t_{k h_{\alpha}} s_{\alpha}$.
Proposition 17.27 The reflection $s_{\alpha, k} \in W$ for all $L_{\alpha, k} \in \mathfrak{Z}$. In fact $s_{\alpha, k}=s_{\alpha-k \delta \delta}$.
Proof. The reflection $s_{\alpha-k \delta}: H_{\mathbb{R}} \rightarrow H_{\mathbb{R}}$ is given by

$$
s_{\alpha-k \delta}(h)=h-(\alpha-k \delta)(h) h_{\alpha-k \delta}
$$

Thus the restriction of $s_{\alpha-k \delta}$ to $H_{\mathbb{R}, 1}$ is given by

$$
s_{\alpha-k \delta}(h)=h-(\alpha(h)-k) h_{\alpha-k \delta}
$$

Since $\delta \in H^{*}$ corresponds to $c \in H$ under our bijection between $H$ and $H^{*}$ we have $h_{\alpha-k \delta}=h_{\alpha}-\frac{2 k}{\langle\alpha, \alpha\rangle} c$. Thus the action of $s_{\alpha-k \delta}$ on $H_{\mathbb{R}, 1} / \mathbb{R} c$ is $s_{\alpha-k \delta}(h)=$ $h-(\alpha(h)-k) h_{\alpha}$ and the action on $H_{\mathbb{R}}^{0}$ is given by the same formula. Thus $s_{\alpha-k \delta}=s_{\alpha, k}$ on $H_{\mathbb{R}}^{0}$. Moreover we know from Theorem 17.17 that $s_{\alpha-k \delta} \in W$ whenever $L_{\alpha, k} \in \mathfrak{Z}$.

We note that $L_{\theta, 1}, L_{\alpha_{1}, 0}, \ldots, L_{\alpha_{l}, 0}$ all lie in $\mathfrak{R}$. For by Proposition $17.18 \theta$ is a short root when $A$ has one of the types $\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}}$.

Definition The connected components of the set $H_{\mathbb{R}}^{0}-\bigcup_{L_{\alpha, k} \in \mathfrak{R}} L_{\alpha, k}$ are called alcoves.

Proposition 17.28 The set

$$
A=\left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \quad \text { for } i=1, \ldots, l, \quad \theta(h)<1\right\}
$$

is an alcove.

Proof. We show $A \cap L_{\alpha, k}=\phi$ for all $L_{\alpha, k} \in \mathcal{R}$. Let $h \in A \cap L_{\alpha, k}$. We may assume $\alpha \in\left(\Phi^{0}\right)^{+}$. Suppose $\theta$ is a long root. Then $0<\alpha(h) \leq \theta(h)<1$ by Proposition 12.9 and so $h$ cannot lie in $L_{\alpha, k}$ for $k \in \mathbb{Z}$. So suppose $\theta$ is a short root. If $\alpha$ is a short root we again have $0<\alpha(h) \leq \theta(h)<1$, so $h$ cannot lie in $L_{\alpha, k}$ with $k \in \mathbb{Z}$. Thus suppose $\alpha$ is a long root. Then $\alpha(h) \leq \theta_{l}(h)$ where $\theta_{l}$ is the highest root of $\Phi^{0}$. Let $\theta_{l}=\sum_{i=1}^{l} b_{i} \alpha_{i}$. Then we have

$$
\begin{aligned}
\theta_{l} & =\sum_{i=1}^{l} b_{i} \alpha_{i} \text { is the highest root of } \Phi^{0} \\
\theta & =\sum_{i=1}^{l} a_{i} \alpha_{i} \text { is the highest short root of } \Phi^{0}
\end{aligned}
$$

By considering the coroot of the highest root, or by a case-by-case check, one may show

$$
b_{i}= \begin{cases}a_{i} & \text { if } \alpha_{i} \text { is long } \\ p a_{i} & \text { if } \alpha_{i} \text { is short }\end{cases}
$$

In particular $b_{i} \leq p a_{i}$ for all $i$. Hence

$$
0<\alpha(h) \leq \theta_{l}(h) \leq p \theta(h)<p
$$

Thus $h$ cannot lie in $L_{\alpha, k}$ with $k \in \mathbb{Z}$ divisible by $p$.
Thus $A$ lies in an alcove. But

$$
\bar{A}=A \cup L_{\alpha_{1}, 0} \cup \cdots \cup L_{\alpha_{l}, 0} \cup L_{\theta, 1}
$$

This shows that $A$ cannot be properly contained in an alcove, since $L_{\alpha_{1}, 0}, \ldots, L_{\alpha_{l}, 0}, L_{\theta, 1}$ lie in $\mathfrak{Z}$. Thus $A$ is an alcove.

Let $\mathfrak{A}$ be the set of alcoves. We show that $W$ acts on $\mathfrak{A}$. Since $W$ is generated by $s_{1}, \ldots, s_{l}, s_{\theta}$ it is sufficient to prove the following lemma.

Lemma 17.29 (i) $s_{i}\left(L_{\alpha, k}\right)=L_{s_{i}(\alpha), k}$ for $i=1, \ldots, l$.
(ii) $s_{\theta}\left(L_{\alpha, k}\right)=L_{s_{0}(\alpha), k+\alpha\left(h_{\theta}\right)}$. Also if $L_{\alpha, k} \in \mathfrak{Z}$ then $L_{s_{0}(\alpha), k+\alpha\left(h_{\theta}\right)} \in \mathfrak{Z}$.

Proof. (i) Let $h \in H_{\mathbb{R}}^{0}$. Then $h \in L_{\alpha, k}$ if and only if $\alpha(h)=k$, and this is equivalent to $\left(s_{i}(\alpha)\right)\left(s_{i}(h)\right)=k$, that is $s_{i}(h) \in L_{s_{i}(\alpha), k}$. Thus $s_{i}\left(L_{\alpha, k}\right)=L_{s_{i}(\alpha), k}$. (ii) $s_{\theta}(h)=s_{0} t_{h_{\theta}}(h)=s_{0}\left(h+h_{\theta}\right)=s_{0}(h)+s_{0}\left(h_{\theta}\right)$. Thus $\alpha(h)=k$ if and only if $\left(s_{0}(\alpha)\right)\left(s_{0}(h)\right)=k$, that is $\left(s_{0}(\alpha)\right)\left(s_{0}(h)+s_{0}\left(h_{\theta}\right)\right)=k+\alpha\left(h_{\theta}\right)$. It follows that $s_{\theta}\left(L_{\alpha, k}\right)=L_{s_{0}(\alpha), k+\alpha\left(h_{\theta}\right)}$.
Now suppose $L_{\alpha, k} \in \mathbb{R}$. Then $k$ is divisible by $p$ if $\alpha$ is a long root and $A \in\left\{\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}}\right\}$. If we are not in this special case then $L_{s_{0}(\alpha), k+\alpha\left(h_{\theta}\right)} \in \mathbb{Z}$ since $\alpha\left(h_{\theta}\right) \in \mathbb{Z}$. So suppose $A$ is one of the above four possibilities and $\alpha$ is a long root. We know $p$ divides $k$ and must show $p$ divides $\alpha\left(h_{\theta}\right)$. Now $h_{\theta}=\frac{1}{a_{0}}\left(c-h_{0}\right)$. We have $a_{0}=1$ in the given cases and $\alpha(c)=0$, thus $\alpha\left(h_{\theta}\right)=-\alpha\left(h_{0}\right)$. Let $\alpha=\sum_{i=1}^{l} k_{i} \alpha_{i}$. Then $\alpha\left(h_{0}\right)=\sum_{i=1}^{l} k_{i} \alpha_{i}\left(h_{0}\right)=\sum_{i=1}^{l} A_{0 i} k_{i}$. There is precisely one $i \in\{1, \ldots, l\}$ with $A_{0 i} \neq 0$. For this $i, A_{0 i}=-2$ in type $\tilde{C}_{l}^{\mathrm{t}}$ and $A_{0 i}=-1$ in the other cases. In the latter cases $\alpha_{i}$ is a short root. Thus

$$
k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{\langle\alpha, \alpha\rangle}=\frac{k_{i}}{p} \in \mathbb{Z}
$$

This shows that $p$ divides $\sum_{i=1}^{l} A_{0 i} k_{i}$, and so $p$ divides $\alpha\left(h_{\theta}\right)$ in all cases. Thus $L_{s_{0}(\alpha), k+\alpha\left(h_{\theta}\right)} \in \mathfrak{R}$.

Corollary 17.30 If $w \in W, A^{\prime} \in \mathfrak{A}$ then $w\left(A^{\prime}\right) \in \mathfrak{A}$.
Proof. This follows from the definition of alcoves, together with the fact that the elements of $W$ permute the affine hyperplanes in $\mathfrak{R}$.

We define $L_{i}=L_{\alpha_{i}, 0}$ for $i=1, \ldots, l$ and $L_{0}=L_{\theta, 1}$. Thus $L_{0}, L_{1}, \ldots, L_{l}$ are the walls bounding the alcove $A$ and $s_{0}, s_{1}, \ldots, s_{l}$ are the reflections in $L_{0}, L_{1}, \ldots, L_{l}$ respectively.

Given $w \in W$ we say that $L_{i}$ separates the alcoves $A$ and $w(A)$ if these alcoves lie on opposite sides of $L_{i}$.

Lemma 17.31 $L_{i}$ separates $A$ and $w(A)$ if and only if $l(w)=l\left(s_{i} w\right)+1$.
Proof. First suppose $w^{\prime} \in W$ has the property that $w^{\prime}(A)$ lies on the same side of $L_{i}$ as $A$ but $w^{\prime} s_{j}(A)$ lies on the opposite side of $L_{i}$ to $A$. Then $w^{\prime}(A), w^{\prime} s_{j}(A)$ lie on opposite sides of $L_{i}$ so $A, s_{j}(A)$ lie on opposite sides of $w^{\prime-1}\left(L_{i}\right)$. This implies $w^{\prime-1}\left(L_{i}\right)=L_{j}$ so $L_{i}=w^{\prime}\left(L_{j}\right)$. Hence $s_{i}=w^{\prime} s_{j} w^{\prime-1}$ and $w^{\prime} s_{j}=s_{i} w^{\prime}$.

Now suppose $w \in W$ is such that $w(A)$ is on the opposite side of $L_{i}$ to $A$. Let $w=s_{i_{1}} \ldots s_{i_{r}}$ be a reduced expression for $w$. Then there exists $q \geq 1$ such that $s_{i_{1}} \ldots s_{i_{q-1}}(A)$ lies on the same side of $L_{i}$ as $A$ but $s_{i_{1}} \ldots s_{i_{q}}(A)$ lies on the opposite side of $L_{i}$. Then we have

$$
s_{i_{1}} \ldots s_{i_{q-1}} s_{i_{q}}=s_{i} s_{i_{1}} \ldots s_{i_{q-1}}
$$

as above. Hence

$$
s_{i} w=s_{i} s_{i_{1}} \ldots s_{i_{r}}=s_{i_{1}} \ldots s_{i_{q-1}} s_{i_{q+1}} \ldots s_{i_{r}}
$$

and so $l\left(s_{i} w\right)<l(w)$.
If $w(A)$ is on the same side of $L_{i}$ as $A$ then $s_{i} w(A)$ is on the opposite side. Hence $l\left(s_{i} \cdot s_{i} w\right)<l\left(s_{i} w\right)$, that is $l\left(s_{i} w\right)>l(w)$.

Theorem 17.32 The map $w \rightarrow w(A)$ is a bijection between the elements of the affine Weyl group $W$ and the set $\mathfrak{A}$ of alcoves.

Proof. Given any alcove $A^{\prime} \in \mathfrak{A}$ we can find a sequence of alcoves

$$
A=A_{1}, A_{2}, \ldots, A_{r}=A^{\prime}
$$

such that $A_{i}$ is obtained from $A_{i-1}$ by reflection in a common wall. Such reflections lie in $W$ by Proposition 17.27. Hence $A^{\prime}=w(A)$ for some $w \in W$. Thus the map $w \rightarrow w(A)$ is surjective.

Next suppose $w(A)=w^{\prime}(A)$. Then $w^{-1} w(A)=A$. We show $w^{\prime-1} w=1$. If this is not so then

$$
w^{\prime-1} w=s_{i} w^{\prime \prime} \quad \text { with } \quad l\left(w^{\prime-1} w\right)=l\left(s_{i} w^{\prime-1} w\right)+1
$$

for some $i$. By Lemma $17.31 L_{i}$ separates $A$ and $w^{\prime-1} w(A)$. This is a contradiction so $w^{\prime-1} w=1$ and $w=w^{\prime}$.

Theorem 17.33 The closure $\bar{A}$ of $A$ is a fundamental region for the action of the affine Weyl group $W$ on $H_{\mathbb{R}}^{0}$, i.e. each $W$-orbit on $H_{\mathbb{R}}^{0}$ intersects $\bar{A}$ in exactly one point.

Proof. Each point in $H_{\mathbb{R}}^{0}$ lies in the closure $\overline{A^{\prime}}$ of some alcove $A^{\prime}$. By Theorem $17.32 A^{\prime}=w(A)$ for some $w \in W$. Thus the $W$-orbit of the given point intersects $\bar{A}$.

Now suppose $x, y \in \bar{A}$ satisfy $y=w(x)$ for $w \in W$. We shall show $x=y$ by induction on $l(w)$. If $l(w)=0$ then $w=1$ so $x=y$. So suppose $l(w)>0$. Then $w=s_{i} w^{\prime}$ with $l\left(s_{i} w\right)<l(w)$. By Lemma $17.31 L_{i}$ separates $A$ and $w(A)$. Thus $\bar{A} \cap w(\bar{A}) \subset L_{i}$. Now $y \in \bar{A} \cap w(\bar{A})$ hence $y \in L_{i}$. Thus $s_{i}(y)=y$. But then $s_{i}(y)=w^{\prime}(x)$ so $y=w^{\prime}(x)$. Since $l\left(w^{\prime}\right)<l(w)$ we deduce $x=y$ by induction.

Remark 17.34 We may also define an action of the affine Weyl group $W$ on $H^{*}$ in a way which is compatible with the bijection $H \rightarrow H^{*}$ determined by the standard invariant form $\langle$,$\rangle on H$. Under this bijection the element $h_{\theta} \in H$ corresponds to $\frac{1}{a_{0}} \theta \in H^{*}$.

For each $\alpha \in\left(H^{0}\right)^{*}$ we may define $t_{\alpha}: H^{*} \rightarrow H^{*}$ by

$$
t_{\alpha}(\lambda)=\lambda+\lambda(c) \alpha-\left(\langle\lambda, \alpha\rangle+\frac{1}{2}\langle\alpha, \alpha\rangle \lambda(c)\right) \delta .
$$

Then we have $s_{0} s_{\theta}=t_{\left(1 / a_{0}\right) \theta}$ on $H^{*}$. Moreover we have $t_{\alpha} t_{\beta}=t_{\alpha+\beta}$ and $w t_{\alpha} w^{-1}=t_{w(\alpha)}$ for $w \in W^{0}$. It follows that we have a semidirect decomposition $W=t\left(M^{*}\right) W^{0}$ where $t\left(M^{*}\right)$ is the set of $t_{\alpha}$ for $\alpha \in M^{*}$ and $M^{*}$ is the sublattice of $\left(H_{\mathbb{R}}^{0}\right)^{*}$ spanned by $w\left(\frac{1}{a_{0}} \theta\right)$ for all $w \in W^{0}$.

The lattice $M^{*}$ is given explicitly as follows.

$$
\begin{aligned}
& M^{*}=\sum_{i=1}^{l} \mathbb{Z} \alpha_{i} \quad \text { for types } \tilde{A}_{l}, \tilde{D}_{l}, \tilde{E}_{6}, \tilde{E}_{7}, \tilde{E}_{8} \\
& M^{*}=\sum_{\alpha_{i} \text { long }} \mathbb{Z} \alpha_{i}+\sum_{\alpha_{i} \text { short }} p \mathbb{Z} \alpha_{i} \quad \text { for types } \tilde{B}_{l}, \tilde{C}_{l}, \tilde{F}_{4}, \tilde{G}_{2} \\
& M^{*}=\sum_{i=1}^{l} \mathbb{Z} \alpha_{i} \quad \text { for types } \tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}} \\
& M^{*}=\sum_{\alpha_{i} \text { long }} \frac{1}{2} \mathbb{Z} \alpha_{i}+\sum_{\alpha_{i} \text { short }} \mathbb{Z} \alpha_{i} \quad \text { for type } \tilde{C}_{l}^{\prime} \\
& M^{*}=\frac{1}{2} \mathbb{Z} \alpha_{1} \quad \text { for type } \tilde{A}_{1}^{\prime} .
\end{aligned}
$$

Now the affine Weyl group $W$ acts on the subset

$$
H_{\mathbb{R}, 1}^{*}=\left\{\lambda \in H_{\mathbb{R}}^{*} ; \quad \lambda(c)=1\right\}
$$

and this induces an action on the orbit space $H_{\mathbb{R}, 1}^{*} / \mathbb{R} \delta$. However, there is a natural bijection between this orbit space and $\left(H_{\mathbb{R}}^{0}\right)^{*}$. This defines a $W$-action on $\left(H_{\mathbb{R}}^{0}\right)^{*}$. The $W_{0}$-action on $\left(H_{\mathbb{R}}^{0}\right)^{*}$ is just as before, and the remaining generator $s_{0}$ of $W$ acts as the reflection in the affine hyperplane

$$
L_{h_{\theta}, 1 / a_{0}}^{*}=\left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{\theta}\right)=1 / a_{0}\right\} .
$$

The element $t_{\alpha}, \alpha \in M^{*}$, acts on $\left(H_{\mathbb{R}}^{0}\right)^{*}$ as translation by $\alpha$, thus $W$ acts on $\left(H_{\mathbb{R}}^{0}\right)^{*}$ as a group of affine transformations.

We may also introduce alcove geometry in $\left(H_{\mathbb{R}}^{0}\right)^{*}$. We define a set $\mathbb{R}^{*}$ of affine hyperplanes in $\left(H_{\mathbb{R}}^{0}\right)^{*}$ as follows.

$$
\mathfrak{R}^{*}=\left\{L_{h_{\alpha}, k}^{*} ; \alpha \in \Phi^{0}, \quad k \text { as below }\right\}
$$

where $L_{h_{\alpha}, k}^{*}=\left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{\alpha}\right)=k\right\}$. The number $k$ runs through the set given as follows.

For types $\tilde{A}_{l}, \tilde{D}_{l}, \tilde{E}_{6}, \tilde{E}_{7}, \tilde{E}_{8} \quad k \in \mathbb{Z}$.
For types $\tilde{B}_{l}, \tilde{C}_{l}, \tilde{F}_{4}, \tilde{G}_{2} \quad k \in \mathbb{Z}$ if $\alpha$ is long $k \in p \mathbb{Z}$ if $\alpha$ is short.
For types $\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}} \quad k \in \mathbb{Z}$.
For type $\tilde{C}_{l}^{\prime} \quad k \in \frac{1}{2} \mathbb{Z} \quad$ if $\alpha$ is long

$$
k \in \mathbb{Z} \quad \text { if } \alpha \text { is short. }
$$

For type $\tilde{A}_{1}^{\prime} \quad k \in \frac{1}{2} \mathbb{Z}$.

Then the elements of the affine Weyl group $W$ permute the set $\mathfrak{R}^{*}$ of affine hyperplanes. The connected components of

$$
H_{\mathbb{R}}^{0}-\bigcup_{L^{*} \in \mathfrak{R}^{*}} L^{*}
$$

are called the alcoves of $H_{\mathbb{R}}^{0}$. The set $A^{*}$ given by

$$
A^{*}=\left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, l, \quad \lambda\left(h_{\theta}\right)<1 / a_{0}\right\}
$$

is an alcove called the fundamental alcove. The group $W$ acts on the alcoves and the map $w \rightarrow w\left(A^{*}\right)$ is a bijective correspondence between elements of $W$ and alcoves. Moreover the closure $\overline{A^{*}}$ is a fundamental region for the $W$-action on $\left(H_{\mathbb{R}}^{0}\right)^{*}$.

We omit the proofs of these facts, which are entirely analogous to the corresponding results for the $W$-action on $H$, or may be deduced from these.

## 18

## Realisations of affine Kac-Moody algebras

### 18.1 Loop algebras and central extensions

Let $A^{0}$ be an indecomposable Cartan matrix of finite type. We have $A^{0}=\left(A_{i j}^{0}\right)$ for $i, j=1, \ldots, l$. Let $L^{0}=L\left(A^{0}\right)$ be the finite dimensional simple Lie algebra with Cartan matrix $A^{0}$. We may construct an $(l+1) \times(l+1)$ affine Cartan matrix $A$ from $A^{0}$ by adding an additional row and column, labelled by 0 , as follows. Let $\theta=\sum_{i=1}^{l} a_{i} \alpha_{i}$ be the highest root of $L^{0}$ and $h_{\theta}=\sum_{i=1}^{l} c_{i} h_{i}$ be the coroot of $\theta$. We then define $A$ by:

$$
\begin{aligned}
& A_{i j}=A_{i j}^{0} \quad \text { if } i, j \in\{1, \ldots, l\} \\
& A_{i 0}=-\sum_{j=1}^{l} a_{j} A_{i j}^{0} \quad \text { if } i \in\{1, \ldots, l\} \\
& A_{0 j}=-\sum_{i=1}^{l} c_{i} A_{i j}^{0} \quad \text { if } j \in\{1, \ldots, l\} \\
& A_{00}=2
\end{aligned}
$$

Proposition 18.1 A is an affine Cartan matrix. The type of $A$ is as follows.

$$
\begin{aligned}
& \text { Type of } A^{0}: A_{l}, B_{l}, C_{l}, D_{l}, E_{6}, E_{7}, E_{8}, F_{4}, G_{2} \\
& \text { Type of } A: \tilde{A}_{l}, \tilde{B}_{l}, \tilde{C}_{l}, \tilde{D}_{l}, \tilde{E}_{6}, \tilde{E}_{7}, \tilde{E}_{8}, \tilde{F}_{4}, \tilde{G}_{2}
\end{aligned}
$$

Proof. We have

$$
A\left(\begin{array}{c}
a_{0} \\
a_{1} \\
\vdots \\
a_{l}
\end{array}\right)=\left(\begin{array}{c}
0 \\
0 \\
\vdots \\
0
\end{array}\right) \quad \text { where } a_{0}=1
$$

For $\sum_{j=0}^{l} A_{i j} a_{j}=A_{i 0}+\sum_{j=1}^{l} A_{i j} a_{j}$. If $i \neq 0$ this is 0 by definition. If $i=0$ we have

$$
\sum_{j=0}^{l} A_{0 j} a_{j}=2+\sum_{j=1}^{l} A_{0 j} a_{j}=2-\sum_{i=1}^{l} \sum_{j=1}^{l} c_{i} A_{i j}^{0} a_{j} .
$$

However, $\quad \sum_{i=1}^{l} \sum_{j=1}^{l} c_{i} A_{i j}^{0} a_{j}=\left(\sum_{j=1}^{l} a_{j} \alpha_{j}\right)\left(\sum_{i=1}^{l} c_{i} h_{i}\right)=\theta\left(h_{\theta}\right)=2$, thus $\sum_{j=0}^{l} A_{0 j} a_{j}=0$.

A similar argument shows that

$$
\left(c_{0} c_{1} \ldots c_{l}\right) A=(00 \ldots 0) \quad \text { where } c_{0}=1
$$

Now $A$ is determined by $A^{0}$ and the relations

$$
A\left(\begin{array}{c}
a_{0} \\
a_{1} \\
\vdots \\
a_{l}
\end{array}\right)=\left(\begin{array}{c}
0 \\
0 \\
\vdots \\
0
\end{array}\right), \quad\left(c_{0} c_{1} \ldots c_{l}\right) A=(0 \ldots 0)
$$

But the affine Cartan matrix $A$ of type $\tilde{L}_{0}$, where $L_{0}$ is $A_{l}, B_{l}, C_{l}, D_{l}, E_{6}, E_{7}$, $E_{8}, F_{4}, G_{2}$ gives $A^{0}$ when row and column 0 are removed, and satisfies the above two relations, by Proposition 17.18 and Lemma 17.19. Thus our given matrix $A$ is the affine Cartan matrix of type $\tilde{L}_{0}$.

Definition An affine Cartan matrix A is of untwisted type if it is one of

$$
\tilde{A}_{l}, \tilde{B}_{l}, \tilde{C}_{l}, \tilde{D}_{l}, \tilde{E}_{6}, \tilde{E}_{7}, \tilde{E}_{8}, \tilde{F}_{4}, \tilde{G}_{2}
$$

Since any affine Cartan matrix $A$ of untwisted type can be constructed as above from a Cartan matrix $A^{0}$ of finite type by the addition of an extra row and column, it seems natural to ask whether the affine Kac-Moody algebra $L(A)$ can be constructed in some way from the finite dimensional simple Lie algebra $L^{0}=L\left(A^{0}\right)$. We shall now describe a method of doing this.

Let $\mathbb{C}\left[t, t^{-1}\right]$ be the ring of Laurent polynomials $\sum_{i \in \mathbb{Z}} \zeta_{i} t^{i}$ for $\zeta_{i} \in \mathbb{C}$ with finitely many $\zeta_{i} \neq 0$. Let

$$
\mathfrak{R}\left(L^{0}\right)=\mathbb{C}\left[t, t^{-1}\right] \otimes_{\mathbb{C}} L^{0} .
$$

Then $\mathfrak{R}\left(L^{0}\right)$ may be made into a Lie algebra in a unique way satisfying

$$
[p \otimes x, q \otimes y]=p q \otimes[x y]
$$

for $p, q \in \mathbb{C}\left[t, t^{-1}\right], \quad x, y \in L^{0}$. This Lie algebra $\mathbb{R}\left(L^{0}\right)$ is called the loop algebra of $L^{0}$.

We now wish to construct a 1-dimensional central extension of $\mathfrak{R}\left(L^{0}\right)$.
Lemma 18.2 Let $L$ be a Lie algebra over $\mathbb{C}$ and $\tilde{L}$ be the set of elements $x+\lambda c$ with $x \in L$ and $\lambda \in \mathbb{C}$. Let $\kappa: L \times L \rightarrow \mathbb{C}$ be a bilinear map satisfying

$$
\begin{aligned}
& \kappa(y, x)=-\kappa(x, y) \quad \text { for } x, y \in L \\
& \kappa([x y], z)+\kappa([y z], x)+\kappa([z x], y)=0 \quad \text { for } x, y, z \in L
\end{aligned}
$$

( $\kappa$ is called a 2-cocycle on L.) Then the Lie multiplication

$$
[x+\lambda c, y+\mu c]=[x y]+\kappa(x, y) c
$$

makes $\tilde{L}$ into a Lie algebra.

Proof. This is elementary. The two relations satisfied by $\kappa$ give anticommutativity and the Jacobi identity on $\tilde{L}$.

We note that $\tilde{L}$ is a 1 -dimensional central extension of $L$, i.e. there is a surjective homomorphism

$$
\theta: \tilde{L} \rightarrow L
$$

given by $\theta(x+\lambda c)=x$, such that $\operatorname{dim}(\operatorname{ker} \theta)=1$ and $\operatorname{ker} \theta$ lies in the centre of $\tilde{L}$.
We apply this idea to construct a 1-dimensional central extension of $\mathbb{Z}\left(L^{0}\right)$ by taking a 2-cocycle on $\mathfrak{R}\left(L^{0}\right)$. Let $\langle$,$\rangle be the invariant bilinear form on L^{0}$ satisfying $\left\langle h_{\theta}, h_{\theta}\right\rangle=2$. Since an invariant bilinear form is determined up to a scalar multiple on $L^{0}$, this condition determines it uniquely. In fact this form on $L^{0}$ is the restriction to $L^{0}$ of the standard invariant form on $L=L(A)$, since for the standard form we have $\langle\theta, \theta\rangle=2$ as in Proposition 17.18, hence

$$
\left\langle h_{\theta}, h_{\theta}\right\rangle=\left\langle\frac{2 \theta}{\langle\theta, \theta\rangle}, \frac{2 \theta}{\langle\theta, \theta\rangle}\right\rangle=\frac{4}{\langle\theta, \theta\rangle}=2
$$

We next define a bilinear form

$$
\langle,\rangle_{t}: \mathbb{R}\left(L^{0}\right) \times \mathbb{R}\left(L^{0}\right) \rightarrow \mathbb{C}\left[t, t^{-1}\right]
$$

by $\langle p \otimes x, q \otimes y\rangle_{t}=p q\langle x, y\rangle$. We define the residue function

$$
\text { Res }: \mathbb{C}\left[t, t^{-1}\right] \rightarrow \mathbb{C}
$$

by $\operatorname{Res}\left(\sum \zeta_{i} t^{i}\right)=\zeta_{-1}$.

Lemma 18.3 The function $\kappa: \mathbb{R}\left(L^{0}\right) \times \mathfrak{R}\left(L^{0}\right) \rightarrow \mathbb{C}$ defined by

$$
\kappa(a, b)=\operatorname{Res}\left\langle\frac{\mathrm{d} a}{\mathrm{~d} t}, b\right\rangle_{t}
$$

is a 2 -cocycle on $\mathbb{R}\left(L^{0}\right)$.

Proof. To show that $\kappa$ is anticommutative it is sufficient to verify that

$$
\kappa\left(t^{i} \otimes x, t^{j} \otimes y\right)=-\kappa\left(t^{j} \otimes y, t^{i} \otimes x\right)
$$

Now

$$
\begin{aligned}
\kappa\left(t^{i} \otimes x, t^{j} \otimes y\right) & =\operatorname{Res}\left\langle i t^{i-1} \otimes x, t^{j} \otimes y\right\rangle_{t} \\
& =\operatorname{Res}\left(i t^{i+j-1}\langle x, y\rangle\right) \\
& = \begin{cases}i\langle x, y\rangle & \text { if } i+j=0 \\
0 & \text { if } i+j \neq 0\end{cases}
\end{aligned}
$$

The anticommutativity follows.
We also need

$$
\begin{aligned}
\kappa\left(\left[t^{i} \otimes x, t^{j} \otimes y\right], t^{k} \otimes z\right) & +\kappa\left(\left[t^{j} \otimes y, t^{k} \otimes z\right], t^{i} \otimes x\right) \\
& +\kappa\left(\left[t^{k} \otimes z, t^{i} \otimes x\right], t^{j} \otimes y\right)=0
\end{aligned}
$$

Now we have

$$
\begin{aligned}
\kappa\left(\left[t^{i}\right.\right. & \left.\left.\otimes x, t^{j} \otimes y\right], t^{k} \otimes z\right)=\kappa\left(t^{i+j} \otimes[x y], t^{k} \otimes z\right) \\
& =\operatorname{Res}\left\langle(i+j) t^{i+j-1} \otimes[x y], t^{k} \otimes z\right\rangle_{t} \\
& =\operatorname{Res}\left((i+j) t^{i+j+k-1}\langle[x y], z\rangle\right) \\
\quad & = \begin{cases}(i+j)\langle[x y], z\rangle & \text { if } i+j+k=0 \\
0 & \text { if } i+j+k \neq 0\end{cases}
\end{aligned}
$$

If $i+j+k \neq 0$ the required property is clear. If $i+j+k=0$ the required sum is

$$
\begin{aligned}
& -k\langle[x y], z\rangle-i\langle[y z], x\rangle-j\langle[z x], y\rangle \\
& \quad=-k\langle[x y], z\rangle-i\langle[x y], z\rangle-j\langle[x y], z\rangle \\
& \quad=0
\end{aligned}
$$

since the form is symmetric and invariant.

We may therefore construct the 1 -dimensional central extension $\tilde{\mathscr{R}}\left(L^{0}\right)$ of $\mathfrak{R}\left(L^{0}\right)$ given by

$$
\tilde{\mathfrak{R}}\left(L^{0}\right)=\mathfrak{R}\left(L^{0}\right) \oplus \mathbb{C} c
$$

whose Lie multiplication is given by

$$
[a+\lambda c, b+\mu c]=[a, b]_{0}+\kappa(a, b) c
$$

where $a, b \in \mathfrak{R}\left(L^{0}\right)$ and $[a, b]_{\tilde{Q}}$ is the Lie product of $a, b$ in $\mathfrak{R}\left(L^{0}\right)$.
We next wish to adjoin to $\tilde{\mathbb{R}}\left(L^{0}\right)$ an element $d$ which acts on $\tilde{\mathbb{R}}\left(L^{0}\right)$ as a derivation.

Lemma 18.4 The map $\Delta: \tilde{\mathfrak{R}}\left(L^{0}\right) \rightarrow \tilde{\mathfrak{R}}\left(L^{0}\right)$ given by $\Delta(a+\lambda c)=t \frac{\mathrm{~d} a}{\mathrm{~d} t}$ for $a \in$ $\mathfrak{R}\left(L^{0}\right), \lambda \in \mathbb{C}$, is a derivation.

Proof. Since $[a+\lambda c, b+\mu c]=[a, b]_{0}+\kappa(a, b) c$ we must show that

$$
t \frac{\mathrm{~d}}{\mathrm{dt}}[a, b]_{0}=\left[t \frac{\mathrm{~d} a}{\mathrm{~d} t}, b+\mu c\right]+\left[a+\lambda c, t \frac{\mathrm{~d} b}{\mathrm{~d} t}\right]
$$

that is

$$
\begin{aligned}
t\left[\frac{\mathrm{~d} a}{\mathrm{~d} t}, b\right]_{0}+t\left[a, \frac{\mathrm{~d} b}{\mathrm{~d} t}\right]_{0}= & t\left[\frac{\mathrm{~d} a}{\mathrm{~d} t}, b\right]_{0}+\kappa\left(t \frac{\mathrm{~d} a}{\mathrm{~d} t}, b\right) c+\left[a, t \frac{\mathrm{~d} b}{\mathrm{~d} t}\right]_{0} \\
& +\kappa\left(a, t \frac{\mathrm{~d} b}{\mathrm{~d} t}\right) c
\end{aligned}
$$

that is $\kappa\left(t \frac{\mathrm{~d} a}{\mathrm{~d} t}, b\right)+\kappa\left(a, t \frac{\mathrm{~d} b}{\mathrm{~d} t}\right)=0$. It is sufficient to prove this when $a=p \otimes x$, $b=q \otimes y$ with $p, q \in \mathbb{C}\left[t, t^{-1}\right], \quad x, y \in L^{0}$. Then

$$
\begin{aligned}
& \kappa\left(t \frac{\mathrm{~d} a}{\mathrm{~d} t}, b\right)+\kappa\left(a, t \frac{\mathrm{~d} b}{\mathrm{~d} t}\right)=\kappa\left(t \frac{\mathrm{~d} p}{\mathrm{~d} t} \otimes x, q \otimes y\right)+\kappa\left(p \otimes x, t \frac{\mathrm{~d} q}{\mathrm{~d} t} \otimes y\right) \\
& \quad=\kappa\left(p \otimes x, t \frac{\mathrm{~d} q}{\mathrm{~d} t} \otimes y\right)-\kappa\left(q \otimes y, t \frac{\mathrm{~d} p}{\mathrm{~d} t} \otimes x\right) \\
& \quad=\operatorname{Res}\left\langle\frac{\mathrm{d} p}{\mathrm{~d} t} \otimes x, t \frac{\mathrm{~d} q}{\mathrm{~d} t} \otimes y\right\rangle_{t}-\operatorname{Res}\left\langle\frac{\mathrm{d} q}{\mathrm{~d} t} \otimes y, t \frac{\mathrm{~d} p}{\mathrm{~d} t} \otimes x\right\rangle_{t} \\
& \quad=\operatorname{Res}\left(t \frac{\mathrm{~d} p}{\mathrm{~d} t} \frac{\mathrm{~d} q}{\mathrm{~d} t}\langle x, y\rangle\right)-\operatorname{Res}\left(t \frac{\mathrm{~d} p}{\mathrm{~d} t} \frac{\mathrm{~d} q}{\mathrm{~d} t}\langle x, y\rangle\right) \\
& \quad=0
\end{aligned}
$$

We now define $\hat{\mathfrak{Z}}\left(L^{0}\right)$ by

$$
\hat{\mathfrak{Z}}\left(L^{0}\right)=\tilde{\mathfrak{R}}\left(L^{0}\right) \oplus \mathbb{C} d
$$

and make $\hat{\mathfrak{R}}\left(L^{0}\right)$ into a Lie algebra by defining the Lie product as

$$
[a+\lambda d, b+\mu d]=[a, b]+\lambda \Delta(b)-\mu \Delta(a)
$$

This is clearly skew-symmetric, and the Jacobi identity follows from the fact that $\Delta$ is a derivation. In particular we have

$$
\begin{aligned}
& {\left[\left(t^{i} \otimes x\right)+\lambda c+\mu d,\left(t^{j} \otimes y\right)+\lambda^{\prime} c+\mu^{\prime} d\right]} \\
& \quad=\left(t^{i+j} \otimes[x y]\right)+\mu j\left(t^{j} \otimes y\right)-\mu^{\prime} i\left(t^{i} \otimes x\right)+\delta_{i,-j} i\langle x, y\rangle c
\end{aligned}
$$

for $x, y \in L^{0}, \quad \lambda, \mu, \lambda^{\prime}, \mu^{\prime} \in \mathbb{C}$.

### 18.2 Realisations of untwisted affine Kac-Moody algebras

We aim to show that $\hat{\mathfrak{R}}\left(L^{0}\right)$ is isomorphic to the affine Kac-Moody algebra $L(A)$. Thus $L(A)$ can be constructed from $L^{0}=L\left(A^{0}\right)$ by the following procedure. First form the loop algebra $\mathfrak{R}\left(L^{0}\right)$. Then form the 1-dimensional central extension $\tilde{\mathfrak{R}}\left(L^{0}\right)$. Finally extend this Lie algebra by a derivation to give $\hat{\mathfrak{R}}\left(L^{0}\right)$.

Theorem 18.5 Let $L^{0}=L\left(A^{0}\right)$ be a finite dimensional simple Lie algebra and let $A$ be the untwisted affine Cartan matrix obtained from $A^{0}$ as in Section 18.1. Then $L(A)$ is isomorphic to $\hat{\mathbb{Z}}\left(L^{0}\right)$.

Proof. We shall define elements $e_{0}, e_{1}, \ldots, e_{l} ; f_{0}, f_{1}, \ldots, f_{l} ; h_{0}, h_{1}, \ldots, h_{l}$ in $\hat{\mathfrak{R}}\left(L^{0}\right)$ with the aim of using Proposition 14.15 to show that our Lie algebra is isomorphic to $L(A)$.

Let $E_{1}, \ldots, E_{l} ; F_{1}, \ldots, F_{l} ; H_{1}, \ldots, H_{l}$ be corresponding generators of $L^{0}$. We define

$$
e_{i}=1 \otimes E_{i}, \quad f_{i}=1 \otimes F_{i}, \quad h_{i}=1 \otimes H_{i}
$$

for $i=1, \ldots, l$. Then $\left[e_{i} f_{i}\right]=h_{i}$ for each $i$. We must also define $e_{0}, f_{0}, h_{0} \in$ $\hat{\mathfrak{Z}}\left(L^{0}\right)$.

We consider the root spaces $L_{\theta}^{0}, L_{-\theta}^{0}$ where $\theta$ is the highest root of $L^{0}$. We have $\operatorname{dim} L_{\theta}^{0}=\operatorname{dim} L_{-\theta}^{0}=1$, and the map $L_{\theta}^{0} \times L_{-\theta}^{0} \rightarrow \mathbb{C}$ given by the invariant bilinear form $\langle$,$\rangle on L^{0}$ defined in Section 18.1 is non-degenerate. Let $\omega^{0}$ be the automorphism of $L^{0}$ satisfying $\omega^{0}\left(E_{i}\right)=-F_{i}, \omega^{0}\left(F_{i}\right)=-E_{i}$. Then $\omega^{0}\left(L_{\theta}^{0}\right)=L_{-\theta}^{0}$. We claim it is possible to choose elements $F_{0} \in L_{\theta}^{0}$, $E_{0} \in L_{-\theta}^{0}$ such that $\omega^{0}\left(F_{0}\right)=-E_{0}$ and $\left\langle F_{0}, E_{0}\right\rangle=1$. First choose any non-zero element $F_{0}^{\prime} \in L_{\theta}^{0}$ and let $E_{0}^{\prime}=-\omega^{0}\left(F_{0}^{\prime}\right)$. Let $\left\langle F_{0}^{\prime}, E_{0}^{\prime}\right\rangle=\xi$. Then we have $\xi \neq 0$.

Now let $F_{0}=\lambda F_{0}^{\prime}$ and $E_{0}=\lambda E_{0}^{\prime}$ for $\lambda \in \mathbb{C}$ with $\lambda \neq 0$. Then we have $E_{0}=$ $-\omega^{0}\left(F_{0}\right)$ and $\left\langle F_{0}, E_{0}\right\rangle=\lambda^{2} \xi$. By a suitable choice of $\lambda \in \mathbb{C}$ we can ensure that $\lambda^{2} \xi=1$.

We now define $e_{0}=t \otimes E_{0}$ and $f_{0}=t^{-1} \otimes F_{0}$. Let $H^{0}$ be the subspace of $L^{0}$ spanned by $h_{1}, \ldots, h_{l}$ and

$$
H=\left(1 \otimes H^{0}\right) \otimes \mathbb{C} c \oplus \mathbb{C} d
$$

We define $h_{0} \in H$ by

$$
h_{0}=\left(1 \otimes\left(-H_{\theta}\right)\right)+c .
$$

Then we have

$$
\left[e_{0} f_{0}\right]=\left[t \otimes E_{0}, t^{-1} \otimes F_{0}\right]=\left(1 \otimes\left[E_{0} F_{0}\right]\right)+\left\langle E_{0}, F_{0}\right\rangle c
$$

But

$$
\left[E_{0} F_{0}\right]=\left\langle E_{0}, F_{0}\right\rangle H_{-\theta}^{\prime}=H_{-\theta}^{\prime}=H_{-\theta}=-H_{\theta}
$$

by Corollary 16.5 , since $\langle\theta, \theta\rangle=2$. Thus

$$
\left[e_{0} f_{0}\right]=\left(1 \otimes\left(-H_{\theta}\right)\right)+c=h_{0}
$$

We also define elements $\alpha_{0}, \alpha_{1}, \ldots, \alpha_{l} \in H^{*}$. We have elements $\alpha_{1}, \ldots, \alpha_{l} \in\left(H^{0}\right)^{*}$ and we extend these to $H^{*}$ by saying that $\alpha_{i}(c)=\alpha_{i}(d)=0$ for $i=1, \ldots, l$. We also define $\theta \in H^{*}$ similarly, saying that $\theta(c)=\theta(d)=0$. Let $\delta \in H^{*}$ be the element defined by

$$
\delta(x)=0 \quad \text { for } x \in H^{0}, \quad \delta(c)=0, \quad \delta(d)=1
$$

We then define $\alpha_{0} \in H^{*}$ by $\alpha_{0}=-\theta+\delta$.
We now show that $\left(H, \Pi, \Pi^{\mathrm{v}}\right)$ is a realisation of $A$ where $\Pi=$ $\left\{\alpha_{0}, \alpha_{1}, \ldots, \alpha_{l}\right\}$ and $\Pi^{\mathrm{v}}=\left\{h_{0}, h_{1}, \ldots, h_{l}\right\} . \Pi$ is linearly independent since $\alpha_{1}, \ldots, \alpha_{l}$ are linearly independent and $\alpha_{0}(d) \neq 0, \alpha_{i}(d)=0$ for $i=1, \ldots, l$. $\Pi^{\mathrm{v}}$ is linearly independent since $h_{1}, \ldots, h_{l}$ are linearly independent and $h_{0}$ involves $c$ whereas $h_{i}$ does not for $i=1, \ldots, l$.

We show that $\alpha_{j}\left(h_{i}\right)=A_{i j}$ for $i, j \in\{0,1, \ldots, l\}$. This is clear if $i \neq 0, j \neq 0$. Also for $i \neq 0$ we have

$$
\alpha_{0}\left(h_{i}\right)=-\theta\left(h_{i}\right)+\delta\left(h_{i}\right)=-\theta\left(h_{i}\right)=-\sum_{j=1}^{l} a_{j} \alpha_{j}\left(h_{i}\right)=-\sum_{j=1}^{l} A_{i j} a_{j}=A_{i 0}
$$

Similarly for $j \neq 0$ we have

$$
\alpha_{j}\left(h_{0}\right)=\alpha_{j}\left(-h_{\theta}+c\right)=-\alpha_{j}\left(h_{\theta}\right)=-\sum_{i=1}^{l} c_{i} \alpha_{j}\left(h_{i}\right)=-\sum_{i=1}^{l} c_{i} A_{i j}=A_{0 j} .
$$

Finally $\alpha_{0}\left(h_{0}\right)=(-\theta+\delta)\left(-h_{\theta}+c\right)=\theta\left(h_{\theta}\right)=2$. Thus $\left(H, \Pi, \Pi^{v}\right)$ is a realisation of $A$.

We next verify the relations

$$
\begin{aligned}
& {\left[e_{i} f_{i}\right]=h_{i}} \\
& {\left[e_{i} f_{j}\right]=0 \quad \text { if } i \neq j} \\
& {\left[x e_{i}\right]=\alpha_{i}(x) e_{i} \quad \text { for } x \in H} \\
& {\left[x f_{i}\right]=-\alpha_{i}(x) f_{i} \quad \text { for } x \in H}
\end{aligned}
$$

where $i=0,1, \ldots, l$. These relations certainly hold when $i \neq 0$ and $j \neq 0$. We have shown above that $\left[e_{0} f_{0}\right]=h_{0}$. For $i \neq 0$ we have

$$
\left[e_{i} f_{0}\right]=\left[1 \otimes E_{i}, t^{-1} \otimes F_{0}\right]=t^{-1} \otimes\left[E_{i} F_{0}\right]=0
$$

since $F_{0} \in L_{\theta}^{0}$ and $\theta$ is the highest root of $L^{0}$. Similarly for $j \neq 0$ we have

$$
\left[e_{0} f_{j}\right]=\left[t \otimes E_{0}, 1 \otimes F_{j}\right]=t \otimes\left[E_{0} F_{j}\right]=0
$$

Now let $x=x_{0}+\lambda c+\mu d \in H$ where $x_{0} \in H^{0}$ and $\lambda, \mu \in \mathbb{C}$. Then

$$
\alpha_{0}(x)=-\theta(x)+\delta(x)=-\theta\left(x_{0}\right)+\mu
$$

since $\theta(c)=\theta(d)=0, \delta\left(x_{0}\right)=\delta(c)=0, \delta(d)=1$. Also

$$
\begin{aligned}
{\left[x e_{0}\right] } & =\left[x_{0}+\lambda c+\mu d, t \otimes E_{0}\right]=\left(t \otimes\left[x_{0} E_{0}\right]\right)+\mu\left(t \otimes E_{0}\right) \\
& =-\theta\left(x_{0}\right)\left(t \otimes E_{0}\right)+\mu\left(t \otimes E_{0}\right) \\
& =\alpha_{0}(x) e_{0}
\end{aligned}
$$

Similarly one shows $\left[x f_{0}\right]=-\alpha_{0}(x) f_{0}$. Thus the required relations are all satisfied.

We show next that $e_{0}, e_{1}, \ldots, e_{l}, f_{0}, f_{1}, \ldots, f_{l}$ and $H$ generate $\hat{\mathfrak{R}}\left(L^{0}\right)$. Let $M$ be the subalgebra of $\hat{\mathfrak{R}}\left(L^{0}\right)$ generated by this subset. Since $E_{1}, \ldots, E_{l}, F_{1}, \ldots, F_{l}$ generate $L^{0}, e_{1}, \ldots, e_{l}, f_{1}, \ldots, f_{l}$ generate $1 \otimes L^{0}$. Thus $1 \otimes L^{0} \subset M$.

Let $I^{0}=\left\{x \in L^{0} ; t \otimes x \in M\right\}$. Since $e_{0}=t \otimes E_{0}$ we have $E_{0} \in I^{0}$ so $I^{0} \neq 0$. Also if $x \in I^{0}, y \in L^{0}$ then $[x y] \in I^{0}$ since

$$
t \otimes[x y]=[t \otimes x, 1 \otimes y] \in M
$$

Thus $I^{0}$ is a non-zero ideal of $L^{0}$. Since $L^{0}$ is simple we have $I^{0}=L^{0}$. Thus $t \otimes x \in M$ for all $x \in L^{0}$. We may now use the relation

$$
\left[t \otimes x, t^{k-1} \otimes y\right]=t^{k} \otimes[x y]
$$

to deduce by induction on $k$ that $t^{k} \otimes x \in M$ for all $x \in L^{0}$ and all $k>0$. In an analogous way, starting with $f_{0}=t^{-1} \otimes F_{0}$ we can show $t^{-k} \otimes x \in M$ for all $x \in L^{0}$ and all $k>0$. Now

$$
\hat{\mathfrak{R}}\left(L^{0}\right)=H+\left(1 \otimes L^{0}\right)+\sum_{k>0}\left(t^{k} \otimes L^{0}\right)+\sum_{k<0}\left(t^{k} \otimes L^{0}\right)
$$

hence $M=\hat{\mathbb{R}}\left(L^{0}\right)$.
It remains to show that $\hat{\mathfrak{R}}\left(L^{0}\right)$ has no non-zero ideal $J$ with $J \cap H=0$. Let

$$
L=\hat{\mathfrak{Z}}\left(L^{0}\right)=H \oplus \sum_{(i, \alpha) \neq(0,0)}\left(t^{i} \otimes\left(L^{0}\right)_{\alpha}\right)
$$

summed over $i \in \mathbb{Z}, \alpha \in\left(H^{0}\right)^{*}$ with $(i, \alpha) \neq(0,0)$. We claim that this is the weight space decomposition of $L$ with respect to $H$. For let $h \in H, x \in\left(L^{0}\right)_{\alpha}$. Then $h=h_{0}+\lambda c+\mu d$ with $h_{0} \in H^{0}, \lambda, \mu \in \mathbb{C}$. Thus

$$
\begin{aligned}
{\left[h, t^{i} \otimes x\right] } & =\left[h_{0}+\lambda c+\mu d, t^{i} \otimes x\right]=\left(t^{i} \otimes\left[h_{0} x\right]\right)+\mu i\left(t^{i} \otimes x\right) \\
& =\left(\alpha\left(h_{0}\right)+\mu i\right)\left(t^{i} \otimes x\right) \\
& =(\alpha(h)+i \delta(h))\left(t^{i} \otimes x\right) \\
& =(\alpha+i \delta)(h)\left(t^{i} \otimes x\right)
\end{aligned}
$$

since $\alpha(h)=\alpha\left(h_{0}\right), \delta(h)=\mu$. Thus $t^{i} \otimes x$ is a weight vector with weight $\alpha+i \delta$. Thus we have

$$
L=L_{0} \oplus \sum_{(\alpha, i) \neq(0,0)} L_{\alpha+i \delta}
$$

where $L_{0}=H$ and $L_{\alpha+i \delta}=t^{i} \otimes\left(L^{0}\right)_{\alpha}$.
Let $J$ be a non-zero ideal of $L$ with $J \cap H=O$. By Lemma 14.12 we have

$$
J=\left(L_{0} \cap J\right) \oplus \sum_{(\alpha, i) \neq(0,0)}\left(L_{\alpha+i \delta} \cap J\right)
$$

Since $L_{0} \cap J=O$ we have $L_{\alpha+i \delta} \cap J \neq O$ for some $(\alpha, i)$. Let $t^{i} \otimes x \in J$ for $x \in\left(L^{0}\right)_{\alpha}$ with $x \neq 0$. Then there exists $y \in\left(L^{0}\right)_{-\alpha}$ with $\langle x, y\rangle \neq 0$. Thus

$$
\left[t^{i} \otimes x, t^{-i} \otimes y\right]=[x y]+i\langle x, y\rangle c
$$

lies in $J \cap H$, and hence

$$
[x y]+i\langle x, y\rangle c=0
$$

Since $[x y] \in H^{0}$ and $\langle x, y\rangle \neq 0$ we must have $i=0$. But this implies $[x y]=0$, whereas we have

$$
[x y]=\langle x, y\rangle h_{\alpha}^{\prime} \neq 0
$$

by Corollary 16.5 . This gives us the required contradiction. Thus $J \cap H=O$ implies $J=O$.

We have now verified all the conditions of Proposition 14.15. We may therefore deduce that $\hat{\mathbb{R}}\left(L^{0}\right)$ is isomorphic to $L(A)$.

We can deduce from this theorem the multiplicities of the imaginary roots of $L(A)$. These multiplicities were not obtained in Chapter 17. We recall from Theorem 16.27 (ii) that the imaginary roots of $L(A)$ have form $k \delta$ where $k \in \mathbb{Z}$ and $k \neq 0$.

Corollary 18.6 Let $A$ be an indecomposable affine GCM of untwisted type. Then the multiplicity of each imaginary root $k \delta, k \neq 0$, is $l=\operatorname{rank} A$.

Proof. We use the realisation $L(A)=\hat{\mathfrak{R}}\left(L^{0}\right)$. The weight space decomposition of $\hat{\mathfrak{Z}}\left(L^{0}\right)$ shows that the root space for the root $k \delta$ is $t^{k} \otimes H^{0}$. The multiplicity of $k \delta$ is the dimension of this root space, which is $\operatorname{dim} H^{0}=l$.

We now make some comments on the isomorphism between $L(A)$ and $\hat{\mathfrak{Z}}\left(L^{0}\right)$ which we have obtained. Firstly the standard invariant form on $L(A)$ maps under this isomorphism to the form on $\hat{\mathfrak{R}}\left(L^{0}\right)$ given as follows:

$$
\begin{aligned}
& \left\langle t^{i} \otimes x, t^{j} \otimes y\right\rangle=0 \quad \text { if } j \neq-i, \quad \text { for } x, y \in L^{0} \\
& \left\langle t^{i} \otimes x, t^{-i} \otimes y\right\rangle=\langle x, y\rangle \\
& \left\langle t^{i} \otimes x, c\right\rangle=0 \\
& \left\langle t^{i} \otimes x, d\right\rangle=0 \\
& \langle c, c\rangle=0 \\
& \langle d, d\rangle=0 \\
& \langle c, d\rangle=1
\end{aligned}
$$

For it is readily checked that the form defined in this way on $\hat{Z}\left(L^{0}\right)$ is invariant. Moreover we also see that the above form on the subspace $\left(1 \otimes H^{0}\right) \oplus$ $\mathbb{C} c \oplus \mathbb{C} d$ of $\hat{\mathbb{R}}\left(L^{0}\right)$ agrees with the standard invariant form on the subspace $H$ of $L(A)$ under our isomorphism between these subspaces. However, the proof of Theorem 16.2 shows that a symmetric invariant bilinear form on $L(A)$ is uniquely determined by its restriction to $H$. Thus the above form on $\hat{\mathfrak{R}}\left(L^{0}\right)$ corresponds to the standard invariant form on $L(A)$.

We also observe that the element $c \in \hat{\mathfrak{Z}}\left(L^{0}\right)$ corresponds to the canonical central element in $L(A)$ under the isomorphism. For we have

$$
h_{0}=\left(1 \otimes-H_{\theta}\right)+c \quad \text { in } \hat{\mathfrak{R}}\left(L^{0}\right)
$$

and hence $\sum_{i=0}^{l} c_{i} h_{i}=c$ in $\hat{\mathfrak{Z}}\left(L^{0}\right)$. It follows that the image of $c$ under the isomorphism is the canonical central element of $L(A)$.

Also, since $\alpha_{0}(d)=1, \alpha_{i}(d)=0$ for $i=1, \ldots, l$ the element $d \in \hat{\mathbb{Z}}\left(L^{0}\right)$ corresponds under the isomorphism to an analogous scaling element $d$ for $L(A)$.

### 18.3 Some graph automorphisms of affine algebras

We now wish to find realisations of the remaining affine Kac-Moody algebras $L(A)$ where $A$ has type $\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}}, \tilde{A}_{1}^{\prime}$ or $\tilde{C}_{l}^{\prime}$. These are called the twisted affine Kac-Moody algebras. We shall obtain realisations for them as fixed point subalgebras of certain automorphisms of untwisted Kac-Moody algebras. Before doing so, however, we consider the graph automorphisms of the untwisted algebras which fix the vertex 0 and therefore arise from graph automorphisms of the corresponding finite dimensional simple Lie algebras. The graph automorphisms of the finite dimensional simple Lie algebras were considered in Section 9.5. We recall from Theorem 9.19 that if $\sigma$ is a graph automorphism of the finite dimensional simple Lie algebra $L(A)$ then $L(A)^{\sigma}$ is isomorphic to the simple Lie algebra $L\left(A^{1}\right)$ where $A^{1}$ is obtained from $A$ as follows.

| $A$ | $:$ | $A_{2 k}$ | $A_{2 k-1}$ | $D_{k+1}$ | $D_{4}$ | $E_{6}$ |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: |
| Order of $\sigma$ | $:$ | 2 | 2 | 2 | 3 | 2 |
| $A^{1}$ | $:$ | $B_{k}$ | $C_{k}$ | $B_{k}$ | $G_{2}$ | $F_{4}$ |

We shall now prove an analogous result to Theorem 9.19 for affine algebras.
Theorem 18.7 Let $A$ be an affine Cartan matrix of type $\tilde{A}_{2 k-1}, \tilde{D}_{k+1}, \tilde{D}_{4}$ or $\tilde{E}_{6}$ and let $\sigma$ be a graph automorphism of the Kac-Moody algebra $L(A)$ which fixes vertex 0 and has order 2, 2, 3, 2 respectively. Let $A^{0}$ be the corresponding finite Cartan matrix and $A^{1}$ be the finite Cartan matrix associated with $A^{0}$ as above. Let $\tilde{A}^{1}$ be the untwisted affine Cartan matrix obtained from $A^{1}$. Then $L(A)^{\sigma}$ is isomorphic to $L\left(\tilde{A}^{1}\right)$.

Specifically we have

$$
\begin{aligned}
L\left(\tilde{A}_{2 k-1}\right)^{\sigma} & \cong L\left(\tilde{C}_{k}\right) \\
L\left(\tilde{D}_{k+1}\right)^{\sigma} & \cong L\left(\tilde{B}_{k}\right) \\
L\left(\tilde{D}_{4}\right)^{\sigma} & \cong L\left(\tilde{G}_{2}\right) \\
L\left(\tilde{E}_{6}\right)^{\sigma} & \cong L\left(\tilde{F}_{4}\right)
\end{aligned}
$$

Proof. The algebra $L\left(A^{0}\right)$ has a Cartan decomposition

$$
L\left(A^{0}\right)=H^{0} \oplus \sum_{\alpha \in \Phi^{0}} L_{\alpha}^{0}
$$

Thus $L(A)$ has a corresponding decomposition

$$
L(A)=H^{0} \oplus \mathbb{C} c \oplus \mathbb{C} d \oplus \sum_{k \neq 0}\left(t^{k} \otimes H^{0}\right)+\sum_{k, \alpha}\left(t^{k} \otimes L_{\alpha}^{0}\right)
$$

Similarly we have decompositions

$$
\begin{aligned}
& L\left(A^{1}\right)=H^{1} \oplus \sum_{\alpha \in \Phi^{1}} L_{\alpha}^{1} \\
& L\left(\tilde{A}^{1}\right)=H^{1} \oplus \mathbb{C} c \oplus \mathbb{C} d \oplus \sum_{k \neq 0}\left(t^{k} \otimes H^{1}\right) \oplus \sum_{k, \alpha}\left(t^{k} \otimes L_{\alpha}^{1}\right) .
\end{aligned}
$$

Consider the graph automorphism $\sigma: L(A) \rightarrow L(A)$. We have

$$
\begin{aligned}
\sigma\left(H^{0}\right) & =H^{0}, \quad \sigma\left(t^{k} \otimes H^{0}\right)=t^{k} \otimes H^{0}, \quad \sigma\left(t^{k} \otimes L_{\alpha}^{0}\right)=t^{k} \otimes L_{\sigma(\alpha)}^{0} \\
\sigma(c) & =c, \quad \sigma(d)=d
\end{aligned}
$$

Hence

$$
L(A)^{\sigma}=\left(H^{0}\right)^{\sigma} \oplus \mathbb{C} c \oplus \mathbb{C} d \oplus \sum_{k \neq 0}\left(t^{k} \otimes\left(H^{0}\right)^{\sigma}\right)+\sum_{k, S}\left(t^{k} \otimes\left(L_{S}^{0}\right)^{\sigma}\right)
$$

where $S$ is an equivalence class of roots in $\Phi^{0}$ and $L_{S}^{0}=\sum_{\alpha \in S} L_{\alpha}^{0}$ (cf. Proposition 9.18). Now the isomorphism $L\left(A^{0}\right)^{\sigma} \rightarrow L\left(A^{1}\right)$ of Theorem 9.19 gives rise to bijective maps

$$
\begin{aligned}
& \left(H^{0}\right)^{\sigma} \rightarrow H^{1} \\
& \left(L_{S}^{0}\right)^{\sigma} \rightarrow L_{\alpha}^{1} \quad \text { where } \alpha \in \Phi^{1} \text { corresponds to } S \\
& t^{k} \otimes\left(H^{0}\right)^{\sigma} \rightarrow t^{k} \otimes H^{1} \\
& t^{k} \otimes\left(L_{S}^{0}\right)^{\sigma} \rightarrow t^{k} \otimes L_{\alpha}^{1} .
\end{aligned}
$$

These maps, together with $c \rightarrow c, d \rightarrow d$, determine a bijective map $\phi: L(A)^{\sigma} \rightarrow L\left(\tilde{A}^{1}\right)$. We wish to show this map is an isomorphism.

Under this bijection $\phi$ the subalgebra $\left(H^{0}\right)^{\sigma} \oplus \mathbb{C} c \oplus \mathbb{C} d$ maps to $H^{1} \oplus \mathbb{C} c \oplus$ $\mathbb{C} d$ and both are abelian. The action of $\left(H^{0}\right)^{\sigma}$ on $t^{k} \otimes\left(H^{0}\right)^{\sigma}$ and $t^{k} \otimes\left(L_{S}^{0}\right)^{\sigma}$ agrees with the action of $H^{1}$ on $t^{k} \otimes H^{1}$ and $t^{k} \otimes L_{\alpha}^{1}$ respectively. The element $c$ lies in the centre on both sides. The action of $d$ on $t^{k} \otimes\left(H^{0}\right)^{\sigma}$ and $t^{k} \otimes\left(L_{S}^{0}\right)^{\sigma}$ (i.e. multiplication by $k$ ) agrees with the action of $d$ on $t^{k} \otimes H^{1}$ and $t^{k} \otimes L_{\alpha}^{1}$ respectively. Thus it is sufficient to compare the multiplication of the root spaces on both sides. These multiplications are trivially preserved by $\phi$ unless
we take two roots whose sum is 0 . So suppose $x, y \in\left(H^{0}\right)^{\sigma}$ and $x_{1}, y_{1}$ are the corresponding elements of $H^{1}$. We have

$$
\begin{aligned}
& {\left[t^{k} \otimes x, t^{-k} \otimes y\right]=k\langle x, y\rangle_{0} c} \\
& {\left[t^{k} \otimes x_{1}, t^{-k} \otimes y_{1}\right]=k\left\langle x_{1}, y_{1}\right\rangle_{1} c}
\end{aligned}
$$

where $\langle,\rangle_{0}$ is the standard invariant form on $L\left(A^{0}\right)$ and $\langle,\rangle_{1}$ the standard invariant form on $L\left(A^{1}\right)$.

Also if $x \in\left(L_{S}^{0}\right)^{\sigma}, y \in\left(L_{-S}^{0}\right)^{\sigma}$ and $x_{1}, y_{1}$ are the corresponding elements of $L_{\alpha}^{1}, L_{-\alpha}^{1}$ then we have

$$
\begin{aligned}
& {\left[t^{k} \otimes x, t^{-k} \otimes y\right]=[x y]+k\langle x, y\rangle_{0} c} \\
& {\left[t^{k} \otimes x_{1}, t^{-k} \otimes y_{1}\right]=\left[x_{1} y_{1}\right]+k\left\langle x_{1}, y_{1}\right\rangle_{1} c .}
\end{aligned}
$$

Thus to show that $\phi$ is an isomorphism it is sufficient to show that the isomorphism $L\left(A^{0}\right)^{\sigma} \rightarrow L\left(A^{1}\right)$ preserves the standard invariant form, that is if $x \rightarrow x_{1}, y \rightarrow y_{1}$ then $\langle x, y\rangle_{0}=\left\langle x_{1}, y_{1}\right\rangle_{1}$. Since any two symmetric invariant bilinear forms on a finite dimensional simple Lie algebra are proportional it is sufficient to check this for just one non-zero value. To do this we choose a 1-element orbit $(i)$ of $\sigma$ on $\{1, \ldots, l\}$. Such a 1-element orbit exists in all the cases being considered. Then we have an element $h_{i} \in L\left(A^{0}\right)^{\sigma}$ mapping to an element $h_{i} \in L\left(A^{1}\right)$. We have

$$
\left\langle h_{i}, h_{i}\right\rangle_{0}=2 \quad \text { and } \quad\left\langle h_{i}, h_{i}\right\rangle_{1}=2 d_{i}=2 a_{i} / c_{i}
$$

A glance at the values of $a_{i}, c_{i}$ for $L\left(A^{1}\right)$ for $i$ coming from 1-element orbits of $\sigma$ shows that $a_{i}=c_{i}$ in these cases, so $d_{i}=1$. Hence $\left\langle h_{i}, h_{i}\right\rangle_{0}=\left\langle h_{i}, h_{i}\right\rangle_{1}$ and it follows that the isomorphism $L\left(A^{0}\right)^{\sigma} \rightarrow L\left(A^{1}\right)$ preserves the standard invariant forms. This completes the proof.
Note The reader will have noticed that the case $L\left(\tilde{A}_{2 k}\right)^{\sigma}$ has not been included in this theorem. The above proof breaks down in this case because $\sigma$ has no 1 -element orbit on $\{1, \ldots, l\}$. The diagrams of $A^{0}, A^{1}$ are as shown.


In fact, if we take the $\sigma$-orbit $(k, k+1)$ on $\{1, \ldots, 2 k\}$, then the isomorphism $L\left(A^{0}\right)^{\sigma} \rightarrow L\left(A^{1}\right)$ of Theorem 9.19 maps

$$
2\left(h_{k}+h_{k+1}\right) \in L\left(A^{0}\right)^{\sigma} \quad \text { to } \quad h_{k} \in L\left(A^{1}\right)
$$

We have

$$
\left\langle h_{k}, h_{k}\right\rangle_{0}=2, \quad\left\langle h_{k+1}, h_{k+1}\right\rangle_{0}=2, \quad\left\langle h_{k}, h_{k+1}\right\rangle_{0}=-1
$$

Thus

$$
\left\langle 2\left(h_{k}+h_{k+1}\right), 2\left(h_{k}+h_{k+1}\right)\right\rangle_{0}=8 .
$$

On the other hand

$$
\left\langle h_{k}, h_{k}\right\rangle_{1}=2 \frac{a_{k}}{c_{k}}=2 d_{k}=4
$$

Thus

$$
\left\langle 2\left(h_{k}+h_{k+1}\right), 2\left(h_{k}+h_{k+1}\right)\right\rangle_{0} \neq\left\langle h_{k}, h_{k}\right\rangle_{1}
$$

and so the isomorphism between $L\left(A_{2 k}\right)^{\sigma}$ and $L\left(B_{k}\right)$ does not preserve the standard invariant form. It does not therefore lead to an isomorphism between $L\left(\tilde{A}_{2 k}\right)^{\sigma}$ and $L\left(\tilde{B}_{k}\right)$ in the manner described in Theorem 18.7.

### 18.4 Realisations of twisted affine algebras

In order to obtain realisations of the twisted Kac-Moody algebras $L(A)$ where $A$ has types $\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}}, \tilde{A}_{1}^{\prime}, \tilde{C}_{l}^{\prime}$ we must consider the fixed point subalgebras of so-called twisted graph automorphisms.

Let $L^{0}=L\left(A^{0}\right)$ be a finite dimensional simple Lie algebra and $\sigma: L^{0} \rightarrow L^{0}$ be a graph automorphism of $L^{0}$. Then $\sigma$ extends to a graph automorphism of $\hat{\mathfrak{Z}}\left(L^{0}\right)=\mathfrak{R}\left(L^{0}\right) \oplus \mathbb{C} c \oplus \mathbb{C} d$ given by

$$
\begin{aligned}
& \sigma\left(t^{i} \otimes x\right)=t^{i} \otimes \sigma(x) \quad \text { for } x \in L^{0} \\
& \sigma(c)=c, \quad \sigma(d)=d
\end{aligned}
$$

Suppose $\sigma$ has order $m$ and let $\delta=\mathrm{e}^{2 \pi \mathrm{i} / m}$. Then we may define an automorphism $\tau$ of $\hat{\mathfrak{Z}}\left(L^{0}\right)$ by

$$
\begin{aligned}
& \tau\left(t^{i} \otimes x\right)=\delta^{-i} t^{i} \otimes \sigma(x) \quad \text { for } \quad x \in L^{0} \\
& \tau(c)=c, \quad \tau(d)=d .
\end{aligned}
$$

$\tau$ is called a twisted graph automorphism of $\hat{\mathfrak{R}}\left(L^{0}\right)$, and also has order $m$. In fact $m=2$ or 3 in the cases which can arise. We shall consider the fixed point subalgebras $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$. In order to do so we first obtain more information about the action of $\sigma$ on $L^{0}$.

Proposition 18.8 (i) Let $L^{0}$ be a simple Lie algebra of type $A_{2 l-1}, D_{l+1}$ or $E_{6}$ and $\sigma$ be a graph automorphism of $L^{0}$ of order 2. Let $\left(L^{0}\right)_{-1}$ be the eigenspace of $\sigma$ on $L^{0}$ with eigenvalue -1 . Then

$$
L^{0}=\left(L^{0}\right)^{\sigma} \oplus\left(L^{0}\right)_{-1}
$$

and $\left(L^{0}\right)_{-1}$ is an irreducible $\left(L^{0}\right)^{\sigma}$-module.
(ii) Let $L^{0}$ have type $D_{4}$ and $\sigma$ be a graph automorphism of $L^{0}$ of order 3 .

Let $\left(L^{0}\right)_{\omega},\left(L^{0}\right)_{\omega^{2}}$ be the eigenspaces of $\sigma$ with eigenvalues $\omega, \omega^{2}$ where $\omega=\mathrm{e}^{2 \pi \mathrm{i} / 3}$. Then

$$
L^{0}=\left(L^{0}\right)^{\sigma} \oplus\left(L^{0}\right)_{\omega} \oplus\left(L^{0}\right)_{\omega^{2}}
$$

and $\left(L^{0}\right)_{\omega},\left(L^{0}\right)_{\omega^{2}}$ are both irreducible $\left(L^{0}\right)^{\sigma}$-modules.
Proof. Let $x \in\left(L^{0}\right)^{\sigma}, y \in\left(L^{0}\right)_{\varepsilon}$ where $\varepsilon$ is an eigenvalue of $\sigma$. Then

$$
\sigma[x y]=[\sigma(x), \sigma(y)]=\varepsilon[x y] .
$$

Thus $[x y] \in\left(L^{0}\right)_{\varepsilon}$ and so $\left(L^{0}\right)_{\varepsilon}$ is an $\left(L^{0}\right)^{\sigma}$-module.
Suppose first that $\sigma$ has order 2. Let $(\alpha, \beta)$ be a 2-element orbit of $\sigma$ on $\Phi^{0}$ and $E_{\alpha}, E_{\beta} \in L^{0}$ be root vectors such that $\sigma\left(E_{\alpha}\right)=E_{\beta}$. Then $E_{\alpha}-E_{\beta} \in$ $\left(L^{0}\right)_{-1}$ and the weight spaces of $\left(L^{0}\right)_{-1}$ are spanned by such elements for all 2-element orbits $(\alpha, \beta)$. The roots $\alpha, \beta \in\left(H^{0}\right)^{*}$ have the same restriction to $\left(\left(H^{0}\right)^{\sigma}\right)^{*}$ and $\alpha \mid\left(\left(H^{0}\right)^{\sigma}\right)^{*}$ is the weight of $E_{\alpha}-E_{\beta}$. The highest weight of the $\left(L^{0}\right)^{\sigma}$-module $\left(L^{0}\right)_{-1}$ is obtained from the highest 2-element orbit $(\alpha, \beta)$. Let us choose the labellings

for the Dynkin diagrams of $A_{2 l-1}, D_{l+1}, E_{6}$. Then the highest 2-element orbits are

$$
\begin{aligned}
& \left(\alpha_{1}+\alpha_{2}+\cdots+\alpha_{2 l-2}, \quad \alpha_{2}+\cdots+\alpha_{2 l-2}+\alpha_{2 l-1}\right) \quad \text { for } A_{2 l-1} \\
& \left(\alpha_{1}+\alpha_{2}+\cdots+\alpha_{l-1}+\alpha_{l}, \quad \alpha_{1}+\alpha_{2}+\cdots+\alpha_{l-1}+\alpha_{l+1}\right) \quad \text { for } D_{l+1} \\
& \left(\alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\alpha_{4}+\alpha_{5}+\alpha_{6}, \quad \alpha_{1}+2 \alpha_{2}+\alpha_{3}+\alpha_{4}+\alpha_{5}+2 \alpha_{6}\right) \quad \text { for } E_{6} .
\end{aligned}
$$

In these three cases the subalgebra $\left(L^{0}\right)^{\sigma}$ has type $C_{l}, B_{l}$ or $F_{4}$ respectively by Theorem 9.19. We choose the labellings

for these Dynkin diagrams. Thus the highest weights for the $\left(L^{0}\right)^{\sigma}$-modules $\left(L^{0}\right)_{-1}$ are

$$
\begin{aligned}
& \alpha_{1}+2 \alpha_{2}+\cdots+2 \alpha_{l-1}+\alpha_{l} \quad \text { for } C_{l} \\
& \alpha_{1}+\alpha_{2}+\cdots+\alpha_{l-1}+\alpha_{l} \\
& \alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+2 \alpha_{4} \quad \text { for } B_{l}
\end{aligned}
$$

Using the equation $\alpha_{i}=\sum_{j} A_{j i} \omega_{j}$ we see that these highest weights are $\omega_{2}$ for $C_{l}, \omega_{1}$ for $B_{l}$ and $\omega_{4}$ for $F_{4}$.

Now $\operatorname{dim}\left(L^{0}\right)_{-1}=\operatorname{dim} L^{0}-\operatorname{dim}\left(L^{0}\right)^{\sigma}$ and this is $2 l^{2}-l-1=\binom{2 l}{2}-1$ for $C_{l}, 2 l+1$ for $B_{l}$, and 26 for $F_{4}$. However, we also have

$$
\begin{aligned}
& \operatorname{dim} L\left(\omega_{2}\right)=\binom{2 l}{2}-1 \quad \text { in } C_{l} \\
& \operatorname{dim} L\left(\omega_{1}\right)=2 l+1 \quad \text { in } B_{l} \\
& \operatorname{dim} L\left(\omega_{4}\right)=26 \quad \text { in } F_{4}
\end{aligned}
$$

by Weyl's dimension formula. Thus in each case $\operatorname{dim}\left(L^{0}\right)_{-1}$ is the dimension of the irreducible module with the appropriate highest weight. Thus $\left(L^{0}\right)_{-1}$ is isomorphic to this irreducible $\left(L^{0}\right)^{\sigma}$-module.

Now suppose $\sigma$ has order 3. Then $L^{0}$ has type $D_{4}$. Let $(\alpha, \beta, \gamma)$ be a 3-element orbit of $\sigma$ on $\Phi^{0}$ and $E_{\alpha}, E_{\beta}, E_{\gamma}$ be root vectors such that $\sigma\left(E_{\alpha}\right)=$ $E_{\beta}, \sigma\left(E_{\beta}\right)=E_{\gamma}$. Then we have

$$
\begin{aligned}
& E_{\alpha}+\omega^{2} E_{\beta}+\omega E_{\gamma} \in\left(L^{0}\right)_{\omega} \\
& E_{\alpha}+\omega E_{\beta}+\omega^{2} E_{\gamma} \in\left(L^{0}\right)_{\omega^{2}}
\end{aligned}
$$

where $\omega=\mathrm{e}^{2 \pi \mathrm{i} / 3}$, and the weight spaces of the $\left(L^{0}\right)^{\sigma}-\operatorname{modules}\left(L^{0}\right)_{\omega}$ and $\left(L^{0}\right)_{\omega^{2}}$ are spanned by such vectors for all 3-element orbits. We choose the labelling

for the Dynkin diagram of $D_{4}$. The highest 3-element orbit of $\sigma$ on $\Phi^{0}$ is then

$$
\left(\alpha_{1}+\alpha_{2}+\alpha_{3}, \quad \alpha_{1}+\alpha_{2}+\alpha_{4}, \quad \alpha_{1}+\alpha_{3}+\alpha_{4}\right)
$$

The subalgebra $\left(L^{0}\right)^{\sigma}$ has type $G_{2}$, for which we take the labelling


Thus the highest weights of the $G_{2}$-modules $\left(L^{0}\right)_{\omega}$ and $\left(L^{0}\right)_{\omega^{2}}$ are both $\alpha_{1}+2 \alpha_{2}$. Now in $G_{2}$ we have $\alpha_{1}+2 \alpha_{2}=\omega_{2}$ and $\operatorname{dim} L\left(\omega_{2}\right)=7$. However, we also have

$$
\operatorname{dim}\left(L^{0}\right)_{\omega}=\operatorname{dim}\left(L^{0}\right)_{\omega^{2}}=\frac{1}{2}\left(\operatorname{dim} L^{0}-\operatorname{dim}\left(L^{0}\right)^{\sigma}\right)=7
$$

Thus the $G_{2}$-modules $\left(L^{0}\right)_{\omega}$ and $\left(L^{0}\right)_{\omega^{2}}$ are both irreducible and isomorphic to $L\left(\omega_{2}\right)$.

Theorem 18.9 Let $L^{0}$ be a simple Lie algebra of type $A_{2 l-1}, D_{l+1}, E_{6}$ or $D_{4}$ and let $\sigma$ be a graph automorphism of $L^{0}$ of order 2, 2, 2, 3 respectively. Let $\tau$ be the corresponding twisted graph automorphism of $\hat{\mathfrak{R}}\left(L^{0}\right)$. Then the fixed point subalgebra $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$ is isomorphic to a twisted affine Kac-Moody algebra.

Explicitly we have

$$
\begin{aligned}
& \hat{\mathfrak{}\left(\tilde{A}_{2 l-1}\right)^{\tau} \cong L\left(\tilde{B}_{l}^{\mathrm{t}}\right)} \\
& \hat{\mathfrak{Q}}\left(\tilde{D}_{l+1}\right)^{\tau} \cong L\left(\tilde{C}_{l}^{\mathrm{t}}\right) \\
& \hat{\mathfrak{}}\left(\tilde{E}_{6}\right)^{\tau} \cong L\left(\tilde{F}_{4}^{\mathrm{t}}\right) \\
& \hat{\mathfrak{R}}\left(\tilde{D}_{4}\right)^{\tau} \cong L\left(\tilde{G}_{2}^{\mathrm{t}}\right) .
\end{aligned}
$$

Proof. The method of proof is broadly similar to that of Theorem 18.5 giving the realisations of the untwisted affine Kac-Moody algebras. The basic idea is to show that the given subalgebra of $\tau$-invariant elements satisfies the conditions of Proposition 14.15, and is therefore isomorphic to the appropriate twisted affine Kac-Moody algebra.

We have

$$
\hat{\mathfrak{Q}}\left(L^{0}\right)=\sum_{k \in \mathbb{Z}}\left(t^{k} \otimes L^{0}\right) \oplus \mathbb{C} c \oplus \mathbb{C} d
$$

If $\sigma$ has order 2 we have

$$
\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}=\sum_{k \in \mathbb{Z}}\left(t^{2 k} \otimes\left(L^{0}\right)^{\sigma}\right) \oplus \sum_{k \in \mathbb{Z}}\left(t^{2 k+1} \otimes\left(L^{0}\right)_{-1}\right) \oplus \mathbb{C} c \oplus \mathbb{C} d
$$

whereas if $\sigma$ has order 3

$$
\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}=\sum_{k \in \mathbb{Z}}\left(t^{3 k} \otimes\left(L^{0}\right)^{\sigma}\right) \oplus \sum_{k \in \mathbb{Z}}\left(t^{3 k+1} \otimes\left(L^{0}\right)_{\omega}\right) \oplus \sum_{k \in \mathbb{Z}}\left(t^{3 k+2} \otimes\left(L^{0}\right)_{\omega^{2}}\right)
$$

$\oplus \mathbb{C} c \oplus \mathbb{C} d$.

Let $E_{1}, \ldots, E_{k}, F_{1}, \ldots, F_{k}, H_{1}, \ldots, H_{k}$ be standard generators of $L^{0}$. We wish to define analogous elements

$$
e_{0}, e_{1}, \ldots, e_{l}, \quad f_{0}, f_{1}, \ldots, f_{l}, \quad h_{0}, h_{1}, \ldots, h_{l} \quad \text { in } \hat{\mathbb{R}}\left(L^{0}\right)^{\tau}
$$

We pick a representative $\theta^{0} \in \Phi^{0}$ of the highest 2- or 3-element $\sigma$-orbit on $\Phi^{0}$. Specifically we have

$$
\begin{aligned}
& \theta^{0}=\alpha_{1}+\alpha_{2}+\cdots+\alpha_{2 l-2} \quad \text { in } A_{2 l-1} \\
& \theta^{0}=\alpha_{1}+\alpha_{2}+\cdots+\alpha_{l} \quad \text { in } D_{l+1} \\
& \theta^{0}=\alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\alpha_{4}+\alpha_{5}+\alpha_{6} \quad \text { in } E_{6} .
\end{aligned}
$$

The elements $e_{i}, f_{i}, h_{i}$ are then chosen as follows.
Type $A_{2 l-1}$

$$
\begin{aligned}
& e_{1}=1 \otimes\left(E_{1}+E_{2 l-1}\right), \ldots, e_{l-1}=1 \otimes\left(E_{l-1}+E_{l+1}\right), \quad e_{l}=1 \otimes E_{l} \\
& f_{1}=1 \otimes\left(F_{1}+F_{2 l-1}\right), \ldots, f_{l-1}=1 \otimes\left(F_{l-1}+F_{l+1}\right), \quad f_{l}=1 \otimes F_{l} \\
& h_{1}=1 \otimes\left(H_{1}+H_{2 l-1}\right), \ldots, h_{l-1}=1 \otimes\left(H_{l-1}+H_{l+1}\right), \quad h_{l}=1 \otimes H_{l} \\
& e_{0}=t \otimes\left(F_{\theta^{0}}-F_{\sigma\left(\theta^{0}\right)}\right), \quad f_{0}=t^{-1} \otimes\left(E_{\theta^{0}}-E_{\sigma\left(\theta^{0}\right)}\right) \\
& h_{0}=1 \otimes\left(-H_{\theta^{0}}-H_{\sigma\left(\theta^{0}\right)}\right)+2 c .
\end{aligned}
$$

Type $D_{l+1}$


$$
\begin{aligned}
& e_{1}=1 \otimes E_{1}, \ldots, e_{l-1}=1 \otimes E_{l-1}, \quad e_{l}=1 \otimes\left(E_{l}+E_{l+1}\right) \\
& f_{1}=1 \otimes F_{1}, \ldots, f_{l-1}=1 \otimes F_{l-1}, \quad f_{l}=1 \otimes\left(F_{l}+F_{l+1}\right) \\
& h_{1}=1 \otimes H_{1}, \ldots, h_{l-1}=1 \otimes H_{l-1}, \quad h_{l}=1 \otimes\left(H_{l}+H_{l+1}\right) \\
& e_{0}=t \otimes\left(F_{\theta^{0}}-F_{\sigma\left(\theta^{0}\right)}\right), \quad f_{0}=t^{-1} \otimes\left(E_{\theta^{0}}-E_{\sigma\left(\theta^{0}\right)}\right) \\
& h_{0}=1 \otimes\left(-H_{\theta^{0}}-H_{\sigma\left(\theta^{0}\right)}\right)+2 c .
\end{aligned}
$$

## Type $E_{6}$



$$
\begin{aligned}
& e_{1}=1 \otimes E_{1}, \quad e_{2}=1 \otimes E_{2}, \quad e_{3}=1 \otimes\left(E_{3}+E_{6}\right), \quad e_{4}=1 \otimes\left(E_{4}+E_{5}\right) \\
& f_{1}=1 \otimes F_{1}, \quad f_{2}=1 \otimes F_{2}, \quad f_{3}=1 \otimes\left(F_{3}+F_{6}\right), \quad f_{4}=1 \otimes\left(F_{4}+F_{5}\right) \\
& h_{1}=1 \otimes H_{1}, \quad h_{2}=1 \otimes H_{2}, \quad h_{3}=1 \otimes\left(H_{3}+H_{6}\right), \quad h_{4}=1 \otimes\left(H_{4}+H_{5}\right) \\
& e_{0}=t \otimes\left(F_{\theta^{0}}-F_{\sigma\left(\theta^{0}\right)}\right), \quad f_{0}=t^{-1} \otimes\left(E_{\theta^{0}}-E_{\sigma\left(\theta^{0}\right)}\right) \\
& h_{0}=1 \otimes\left(-H_{\theta^{0}}-H_{\sigma\left(\theta^{0}\right)}\right)+2 c .
\end{aligned}
$$

Type $D_{4}$


$$
\begin{aligned}
& e_{1}=1 \otimes E_{1}, \quad e_{2}=1 \otimes\left(E_{2}+E_{3}+E_{4}\right) \\
& f_{1}=1 \otimes F_{1}, \quad f_{2}=1 \otimes\left(F_{2}+F_{3}+F_{4}\right) \\
& h_{1}=1 \otimes H_{1}, \quad h_{2}=1 \otimes\left(H_{2}+H_{3}+H_{4}\right) \\
& e_{0}=t \otimes\left(F_{\theta^{0}}+\omega^{2} F_{\sigma\left(\theta^{0}\right)}+\omega F_{\sigma^{2}\left(\theta^{0}\right)}\right) \\
& f_{0}=t^{-1} \otimes\left(E_{\theta^{0}}+\omega E_{\sigma\left(\theta^{0}\right)}+\omega^{2} E_{\sigma^{2}\left(\theta^{0}\right)}\right) \\
& h_{0}=1 \otimes\left(-H_{\theta^{0}}-H_{\sigma\left(\theta^{0}\right)}-H_{\sigma^{2}\left(\theta^{0}\right)}\right)+3 c .
\end{aligned}
$$

Let $H \subset \hat{\mathfrak{R}}\left(L^{0}\right)$ be given by $H=\left(1 \otimes H^{0}\right) \oplus \mathbb{C} c \oplus \mathbb{C} d$. We define maps $\alpha_{0}, \alpha_{1}, \ldots, \alpha_{l}: H^{\sigma} \rightarrow \mathbb{C}$. We have roots $\alpha_{i} \in H^{*}$ and such roots in the same $\sigma$-orbit have the same restriction to $H^{\sigma}$. We define $\alpha_{1}, \ldots, \alpha_{l} \in\left(H^{\sigma}\right)^{*}$ to be the restrictions of the corresponding roots in $H^{*}$. We also define $\alpha_{0} \in\left(H^{\sigma}\right)^{*}$ by $\alpha_{0}=-\theta^{0}+\delta$.

Let $\hat{\mathfrak{R}}\left(L^{0}\right) \cong L(A)$ as in Theorem 18.5. We wish to show that $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau} \cong$ $L\left(A^{\prime}\right)$ where $A^{\prime}$ is the affine Cartan matrix of type given below:

$$
\begin{array}{lllll}
A: & \tilde{A}_{2 l-1} & \tilde{D}_{l+1} & \tilde{E}_{6} & \tilde{D}_{4} \\
A^{\prime}: & \tilde{B}_{l}^{\mathrm{t}} & \tilde{C}_{l}^{\mathrm{t}} & \tilde{F}_{4}^{\mathrm{t}} & \tilde{G}_{2}^{\mathrm{t}}
\end{array}
$$

We shall first show that $\left(H^{\sigma}, \Pi, \Pi^{v}\right)$ is a realisation of $A^{\prime}$ where

$$
\Pi^{\mathrm{v}}=\left\{h_{0}, h_{1}, \ldots, h_{l}\right\} \subset H^{\sigma}, \quad \Pi=\left\{\alpha_{0}, \alpha_{1}, \ldots, \alpha_{l}\right\} \subset\left(H^{\sigma}\right)^{*}
$$

We know from Theorem 9.19 that $\alpha_{j}\left(h_{i}\right)=A_{i j}^{\prime}$ for $i, j \in\{1, \ldots, l\}$. This $l \times l$ matrix is non-singular and so $h_{1}, \ldots, h_{l}$ and $\alpha_{1}, \ldots, \alpha_{l}$ are linearly independent. The element $h_{0}$ involves $c$ whereas $h_{1}, \ldots, h_{l}$ do not, thus $h_{0}, h_{1}, \ldots, h_{l}$ are linearly independent. We have $\alpha_{i}(d)=0$ for $i=1, \ldots, l$ but $\alpha_{0}(d) \neq 0$, thus $\alpha_{0}, \alpha_{1}, \ldots, \alpha_{l}$ are linearly independent. We must show that

$$
\begin{array}{ll}
\alpha_{j}\left(h_{0}\right)=A_{0 j}^{\prime} & j=1, \ldots, l \\
\alpha_{0}\left(h_{i}\right)=A_{i 0}^{\prime} & i=1, \ldots, l \\
\alpha_{0}\left(h_{0}\right)=2 &
\end{array}
$$

We recall the integers $a_{i}, c_{i}$ associated with the affine Cartan matrix $A^{\prime}$, which are as follows.

$$
a_{0}, a_{1}, \ldots, a_{l}
$$




We then note that the following significant equations hold in each of the cases being considered:

$$
\begin{aligned}
\sum_{i=0}^{l} a_{i} \alpha_{i} & =\delta \\
\sum_{i=0}^{l} c_{i} h_{i} & =m c \quad \text { where } m \text { is the order of } \sigma .
\end{aligned}
$$

We then have

$$
\alpha_{0}\left(h_{i}\right)=-\sum_{j=1}^{l} a_{j} \alpha_{j}\left(h_{i}\right)=-\sum_{j=1}^{l} A_{i j}^{\prime} a_{j}=A_{i 0}^{\prime} a_{0}=A_{i 0}^{\prime}
$$

for $i=1, \ldots, l$ since $a_{0}=1$ in the cases being considered. Similarly

$$
\alpha_{j}\left(h_{0}\right)=-\sum_{i=1}^{l} c_{i} \alpha_{j}\left(h_{i}\right)=-\sum_{i=1}^{l} c_{i} A_{i j}^{\prime}=c_{0} A_{0 j}^{\prime}=A_{0 j}^{\prime}
$$

for $j=1, \ldots, l$. Also

$$
\alpha_{0}\left(h_{0}\right)=\left(-\theta^{0}+\delta\right)\left(1 \otimes\left(H_{\theta^{0}}-H_{\sigma\left(\theta^{0}\right)}-\cdots\right)+m c\right)=\theta^{0}\left(H_{\theta^{0}}\right)=2
$$

We also note that $A^{\prime}$ is an $(l+1) \times(l+1)$ matrix of rank $l$ and that $\operatorname{dim} H^{\sigma}=$ $l+2$. Thus we have shown that $\left(H^{\sigma}, \Pi, \Pi^{v}\right)$ is a realisation of $A^{\prime}$.

We next verify the relations necessary for applying Proposition 14.15. We first show that

$$
\left[h_{i} e_{j}\right]=A_{i j}^{\prime} e_{j} \quad\left[h_{i} f_{j}\right]=-A_{i j}^{\prime} f_{j}
$$

for $i, j \in\{0,1, \ldots, l\}$. These are known for $i, j \in\{1, \ldots, l\}$ by Theorem 9.19. So we must check

$$
\begin{aligned}
& {\left[h_{0} e_{j}\right]=A_{0 j}^{\prime} e_{j}, \quad\left[h_{0} f_{j}\right]=-A_{0 j}^{\prime} f_{j}} \\
& {\left[h_{i} e_{0}\right]=A_{i 0}^{\prime} e_{0}, \quad\left[h_{i} f_{0}\right]=-A_{i 0}^{\prime} f_{0}} \\
& {\left[h_{0} e_{0}\right]=2 e_{0}, \quad\left[h_{0} f_{0}\right]=-2 f_{0} .}
\end{aligned}
$$

For $j=1, \ldots, l$ we have

$$
\left[h_{0} e_{j}\right]=-\sum_{i=1}^{l} c_{i}\left[h_{i} e_{j}\right]=\left(-\sum_{i=1}^{l} c_{i} A_{i j}^{\prime}\right) e_{j}=A_{0 j}^{\prime} e_{j}
$$

and similarly for $i=1, \ldots, l$ we have $\left[h_{0} f_{j}\right]=-A_{0 j}^{\prime} f_{j}$. Also

$$
\begin{aligned}
{\left[h_{i} e_{0}\right] } & =\left[h_{i}, t \otimes\left(F_{\theta^{0}}+\varepsilon^{-1} F_{\sigma\left(\theta^{0}\right)}+\cdots\right)\right] \\
& =t \otimes-\theta^{0}\left(h_{i}\right)\left(F_{\theta^{0}}+\varepsilon^{-1} F_{\sigma\left(\theta^{0}\right)}+\cdots\right) \\
& =-\theta^{0}\left(h_{i}\right) e_{0}=\left(-\sum_{j=1}^{l} a_{j} \alpha_{j}\left(h_{i}\right)\right) e_{0}=\left(-\sum_{j=1}^{l} A_{i j}^{\prime} a_{j}\right) e_{0} \\
& =A_{i 0}^{\prime} e_{0} .
\end{aligned}
$$

Similarly we have $\left[h_{i} f_{0}\right]=-A_{i 0}^{\prime} f_{0}$. We also have

$$
\begin{aligned}
{\left[h_{0} e_{0}\right] } & =\left[1 \otimes\left(-H_{\theta^{0}}-H_{\sigma\left(\theta^{0}\right)}-\cdots\right), \quad t \otimes\left(F_{\theta^{0}}+\varepsilon^{-1} F_{\sigma\left(\theta^{0}\right)}+\cdots\right)\right] \\
& =2 t \otimes\left(F_{\theta^{0}}+\varepsilon^{-1} F_{\sigma\left(\theta^{0}\right)}+\cdots\right)=2 e_{0}
\end{aligned}
$$

Similarly we have $\left[h_{0} f_{0}\right]=-2 f_{0}$.

Finally we have relations

$$
\begin{aligned}
& {\left[c e_{i}\right]=0=\alpha_{i}(c) e_{i}} \\
& {\left[c f_{i}\right]=0=-\alpha_{i}(c) f_{i}} \\
& {[d=1, \ldots, l} \\
& {\left[d e_{i}\right]=0=\alpha_{i}(d) e_{i} \quad i=1, \ldots, l} \\
& {\left[d f_{i}\right]=0=-\alpha_{i}(d) f_{i} \quad i=1, \ldots, l} \\
& {\left[d e_{0}\right]=e_{0}=\alpha_{0}(d) e_{0}} \\
& {\left[d f_{0}\right]=-f_{0}=-\alpha_{0}(d) f_{0}}
\end{aligned}
$$

Since $H^{\tau}=\left(1 \otimes\left(H^{0}\right)^{\sigma}\right) \oplus \mathbb{C} c \oplus \mathbb{C} d$ we have now verified all relations necessary for applying Proposition 14.15.

We next show that the elements $e_{0}, e_{1}, \ldots, e_{l}, f_{0}, f_{1}, \ldots, f_{l}$ together with $H^{\tau}$ generate $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$. We know that $e_{1}, \ldots, e_{l}, f_{1}, \ldots, f_{l}$ generate $\left(L^{0}\right)^{\sigma}$ by Theorem 9.19. Since

$$
\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}=\sum_{k \in \mathbb{Z}}\left(t^{2 k} \otimes\left(L^{0}\right)^{\sigma}\right) \oplus \sum_{k \in \mathbb{Z}}\left(t^{2 k+1} \otimes\left(L^{0}\right)_{-1}\right) \oplus \mathbb{C} c \oplus \mathbb{C} d
$$

when $\sigma$ has order 2 and
$\hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}=\sum_{k \in \mathbb{Z}}\left(t^{3 k} \otimes\left(L^{0}\right)^{\sigma}\right) \oplus \sum_{k \in \mathbb{Z}}\left(t^{3 k+1} \otimes\left(L^{0}\right)_{\omega}\right) \oplus \sum_{k \in \mathbb{Z}}\left(t^{3 k+2} \otimes\left(L^{0}\right)_{\omega^{2}}\right) \oplus \mathbb{C} c \oplus \mathbb{C} d$
when $\sigma$ has order 3, it is sufficient to show that the subspaces $\left(t^{2 k} \otimes\left(L^{0}\right)^{\sigma}\right)$ for $k \neq 0$ and $t^{2 k+1} \otimes\left(L^{0}\right)_{-1}$ lie in the subalgebra generated by the above elements when $\sigma$ has order 2, and the subspaces $\left(t^{3 k} \otimes\left(L^{0}\right)^{\sigma}\right)$ for $k \neq 0$, $\left(t^{3 k+1} \otimes\left(L^{0}\right)_{\omega}\right),\left(t^{3 k+2} \otimes\left(L^{0}\right)_{\omega^{2}}\right)$ lie in this subalgebra when $\sigma$ has order 3.

Let $M$ be the subalgebra of $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$ generated by $e_{0}, e_{1}, \ldots, e_{l}$, $f_{0}, f_{1}, \ldots, f_{l}, H^{\tau}$. Suppose first that $\sigma$ has order 2 . We have

$$
e_{0}=t \otimes\left(F_{\theta^{0}}-F_{\sigma\left(\theta^{0}\right)}\right) \in M \text { and } F_{\theta^{0}}-F_{\sigma\left(\theta^{0}\right)} \in\left(L^{0}\right)_{-1}
$$

Now if $x \in\left(L^{0}\right)^{\sigma}, y \in\left(L^{0}\right)_{-1}$ then

$$
[1 \otimes x, t \otimes y]=t \otimes[x y] \in t \otimes\left(L^{0}\right)_{-1}
$$

Thus the elements $y \in\left(L^{0}\right)_{-1}$ for which $t \otimes y \in M$ form an $\left(L^{0}\right)^{\sigma}$-submodule of $\left(L^{0}\right)_{-1}$. However, $\left(L^{0}\right)_{-1}$ is an irreducible $\left(L^{0}\right)^{\sigma}$-module by Proposition 18.8. Thus $t \otimes\left(L^{0}\right)_{-1}$ lies in $M$. Now we can find elements $x, y \in\left(L^{0}\right)_{-1}$ such that $[x y] \neq 0$. Then $[t \otimes x, t \otimes y]=t^{2} \otimes[x y]$ is a non-zero element of $M$. However,
the set of $z \in\left(L^{0}\right)^{\sigma}$ such that $t^{2} \otimes z \in M$ is an ideal of $\left(L^{0}\right)^{\sigma}$ and $\left(L^{0}\right)^{\sigma}$ is a simple Lie algebra. Thus $t^{2} \otimes\left(L^{0}\right)^{\sigma}$ lies in $M$. The relations

$$
\begin{array}{lr}
{\left[t^{2} \otimes x, t^{2 k} \otimes y\right]=t^{2 k+2} \otimes[x y]} & x, y \in\left(L^{0}\right)^{\sigma} \\
{\left[t^{2} \otimes x, t^{2 k+1} \otimes y\right]=t^{2 k+3} \otimes[x y]} & x \in\left(L^{0}\right)^{\sigma}, y \in\left(L^{0}\right)_{-1}
\end{array}
$$

can then be used to show by induction on $k$ that $t^{2 k} \otimes\left(L^{0}\right)^{\sigma} \subset M$ and $t^{2 k+1} \otimes\left(L^{0}\right)_{-1} \subset M$ when $k>0$. Starting with $f_{0}$ instead of $e_{0}$ will similarly show this when $k<0$. Thus $M=\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$.

Now suppose that $\sigma$ has order 3 . We have

$$
\begin{aligned}
& e_{0}=t \otimes\left(F_{\theta^{0}}+\omega^{2} F_{\sigma\left(\theta^{0}\right)}+\omega F_{\sigma^{2}\left(\theta^{0}\right)}\right) \\
& F_{\theta^{0}}+\omega^{2} F_{\sigma\left(\theta^{0}\right)}+\omega F_{\sigma^{2}\left(\theta^{0}\right)} \in\left(L^{0}\right)_{\omega} .
\end{aligned}
$$

An argument similar to the above shows that $t \otimes\left(L^{0}\right)_{\omega} \subset M$. Now there exist elements $x, y \in\left(L^{0}\right)_{\omega}$ with $[x y] \neq 0$. Then

$$
[t \otimes x, t \otimes y]=t^{2} \otimes[x y] \in t^{2} \otimes\left(L^{0}\right)_{\omega^{2}}
$$

We can then show as above that $t^{2} \otimes\left(L^{0}\right)_{\omega^{2}} \subset M$. There exist elements $x \in\left(L^{0}\right)_{\omega}, y \in\left(L^{0}\right)_{\omega^{2}}$ with $[x y] \neq 0$. Then

$$
\left[t \otimes x, t^{2} \otimes y\right]=t^{3} \otimes[x y] \in t^{3} \otimes\left(L^{0}\right)^{\sigma}
$$

We then see as above that $t^{3} \otimes\left(L^{0}\right)^{\sigma} \subset M$. Induction on $k$ can then be used to see that the subspaces

$$
t^{3 k} \otimes\left(L^{0}\right)^{\sigma}, \quad t^{3 k+1} \otimes\left(L^{0}\right)_{\omega}, \quad t^{3 k+2} \otimes\left(L^{0}\right)_{\omega^{2}}
$$

all lie in $M$ when $k>0$. A similar result is obtained when $k<0$ starting with $f_{0}$ instead of $e_{0}$. Thus $M=\hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}$. Hence in all cases the elements $e_{0}, e_{1}, \ldots, e_{l}, f_{0}, f_{1}, \ldots, f_{l}$ and $H^{\tau}$ generate $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$.

Finally we must show that $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$ has no non-zero ideal $J$ with $J \cap H^{\tau}=O$. To see this we decompose $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$ into root spaces with respect to $H^{\tau}$. We first suppose $\sigma$ has order 2. For each 1-element orbit $(\alpha)$ of $\sigma$ on $\Phi^{0}$ we choose $E_{\alpha} \in L_{\alpha}^{0}$. We showed in the proof of Theorem 9.19 that $\sigma\left(E_{\alpha}\right)=E_{\alpha}$. For each 2-element orbit $(\alpha, \beta)$ of $\sigma$ on $\Phi^{0}$ we choose $E_{\alpha} \in L_{\alpha}^{0}, E_{\beta} \in L_{\beta}^{0}$ such that $\sigma\left(E_{\alpha}\right)=E_{\beta}$. Then $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$ is the direct sum of $H^{\tau}$ and the following weight spaces.
$t^{2 k} \otimes\left(H^{0}\right)^{\sigma} \quad$ with weight $2 k \delta$
$t^{2 k+1} \otimes\left(H^{0}\right)_{-1} \quad$ with weight $(2 k+1) \delta$
$t^{2 k} \otimes \mathbb{C} E_{\alpha} \quad$ with weight $\alpha+2 k \delta$ where $(\alpha)$ is a 1-element orbit
$t^{2 k} \otimes \mathbb{C}\left(E_{\alpha}+E_{\beta}\right) \quad$ with weight $\alpha+2 k \delta$ where $(\alpha, \beta)$ is a 2-element orbit
$t^{2 k+1} \otimes \mathbb{C}\left(E_{\alpha}-E_{\beta}\right) \quad$ with weight $\alpha+(2 k+1) \delta$ where $(\alpha, \beta)$ is a 2-element orbit.

By Lemma 14.12 the ideal $J$ is the direct sum of its intersections with these weight spaces. Thus $J$ has non-zero intersection with one of these weight spaces. Taking a non-zero element $x$ in this weight space and in $J$ we can find an element $y$ in the negative weight space such that $[x y]$ is a non-zero element of $H^{\tau}$, and this contradicts $J \cap H^{\tau}=O$.

When $\sigma$ has order 3 a similar argument can be applied. This time the weight spaces are
$t^{3 k} \otimes\left(H^{0}\right)^{\sigma} \quad$ with weight $3 k \delta$
$t^{3 k+1} \otimes\left(H^{0}\right)_{\omega} \quad$ with weight $(3 k+1) \delta$
$t^{3 k+2} \otimes\left(H^{0}\right)_{\omega^{2}} \quad$ with weight $(3 k+2) \delta$
$t^{3 k} \otimes \mathbb{C} E_{\alpha} \quad$ with weight $\alpha+3 k \delta$ where $(\alpha)$ is a 1-element orbit
$t^{3 k} \otimes \mathbb{C}\left(E_{\alpha}+E_{\beta}+E_{\gamma}\right) \quad$ with weight $\alpha+3 k \delta$ where $(\alpha, \beta, \gamma)$ is a
3-element orbit
$t^{3 k+1} \otimes \mathbb{C}\left(E_{\alpha}+\omega^{2} E_{\beta}+\omega E_{\gamma}\right) \quad$ with weight $\alpha+(3 k+1) \delta$ where $(\alpha, \beta, \gamma)$ is a 3-element orbit
$t^{3 k+2} \otimes \mathbb{C}\left(E_{\alpha}+\omega E_{\beta}+\omega^{2} E_{\gamma}\right) \quad$ with weight $\alpha+(3 k+2) \delta$ where $(\alpha, \beta, \gamma)$ is a 3-element orbit.

Any non-zero ideal $J$ with $J \cap H^{\tau}=O$ must have non- zero intersection with one of these weight spaces. We can then multiply it by an element of the negative weight space to give a non-zero element of $H^{\tau}$, and this contradicts $J \cap H^{\tau}=O$. Thus we deduce that $J=O$.

We have now verified all the conditions of Proposition 14.15 and can conclude that $\hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}$ is isomorphic to $L\left(A^{\prime}\right)$.

As a corollary we obtain the multiplicities of the imaginary roots of the affine Kac-Moody algebras of types $\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}}$.

Corollary 18.10 The multiplicities of the imaginary roots are as follows.
Type $\tilde{B}_{l}^{\mathrm{t}}$. The roots $2 k \delta(k \neq 0)$ have multiplicity $l$ and the roots $(2 k+1) \delta$ have multiplicity $l-1$.
Type $\tilde{C}_{l}^{\mathrm{t}}$ The roots $2 k \delta(k \neq 0)$ have multiplicity $l$ and the roots $(2 k+1) \delta$ have multiplicity 1 .
Type $\tilde{F}_{4}^{t}$ The roots $2 k \delta(k \neq 0)$ have multiplicity 4 and the roots $(2 k+1) \delta$ have multiplicity 2 .
Type $\tilde{G}_{2}^{t}$ The roots $3 k \delta(k \neq 0)$ have multiplicity 2 and the roots $(3 k+1) \delta$ and $(3 k+2) \delta$ have multiplicity 1 .

Proof. For types $\tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}$ Theorem 18.9 shows that the multiplicity of $2 k \delta(k \neq 0)$ is $\operatorname{dim}\left(H^{0}\right)^{\sigma}$, which is equal to $l$. The multiplicity of $(2 k+1) \delta$ is $\operatorname{dim}\left(H^{0}\right)_{-1}$, which is $l-1,1,2$ in the three cases respectively.

For type $G_{2}^{t}$ the multiplicity of $3 k \delta(k \neq 0)$ is $\operatorname{dim}\left(H^{0}\right)^{\sigma}=2$, and the multiplicities of $(3 k+1) \delta$ and $(3 k+2) \delta$ are $\operatorname{dim}\left(H^{0}\right)_{\omega}=\operatorname{dim}\left(H^{0}\right)_{\omega^{2}}=1$.

We now make some comments on the isomorphism between $L\left(A^{\prime}\right)$ and $\hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}$ which we have obtained. The standard invariant form on $L\left(A^{\prime}\right)$ does not map under this isomorphism to the restriction of the standard invariant form on $\hat{\mathfrak{R}}\left(L^{0}\right)$. For let (i) be a 1 -element $\sigma$-orbit on $\{1, \ldots, l\}$. Such an orbit exists in each of the cases. Then $h_{i} \in L\left(A^{\prime}\right)$ corresponds to $1 \otimes H_{i} \in \hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}$. We have

$$
\begin{aligned}
& \left\langle 1 \otimes H_{i}, 1 \otimes H_{i}\right\rangle=2 \\
& \left\langle h_{i}, h_{i}\right\rangle^{\prime}=2 a_{i} / c_{i} .
\end{aligned}
$$

However, we may check that $c_{i}=m a_{i}$ for all 1 -element $\sigma$-orbits, hence

$$
\left\langle h_{i}, h_{i}\right\rangle^{\prime}=2 / m
$$

where $m$ is the order of $\sigma$. Thus

$$
\left\langle 1 \otimes H_{i}, 1 \otimes H_{i}\right\rangle=m\left\langle h_{i}, h_{i}\right\rangle^{\prime} .
$$

Also the isomorphism does not map the element $c \in \hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}$ to the canonical central element $c^{\prime} \in L\left(A^{\prime}\right)$. We noted in the proof of Theorem 18.9 that

$$
\sum_{i=0}^{l} c_{i} h_{i}=m c \quad \text { in } \hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}
$$

and hence our isomorphism maps $m c$ to $c^{\prime}$.
However, the scaling element $d \in \hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}$ maps to a scaling element $d^{\prime} \in$ $L\left(A^{\prime}\right)$. For we have $\alpha_{0}=-\theta^{0}+\delta$ and so $\alpha_{0}(d)=\delta(d)=1$. Also $\alpha_{i}(d)=0$ for $i=1, \ldots, l$.

We also see that, if $x, y \in \hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$ map to $x^{\prime}, y^{\prime} \in L\left(A^{\prime}\right)$, then

$$
\langle x, y\rangle=m\left\langle x^{\prime}, y^{\prime}\right\rangle^{\prime}
$$

where $\langle$,$\rangle and \langle,\rangle^{\prime}$ are the standard invariant forms. This is true for $x^{\prime}, y^{\prime} \in$ $L\left(A^{\prime}\right)^{\prime}$ since any two symmetric invariant forms are proportional on $L\left(A^{\prime}\right)^{\prime}$. However, it is also true when $y=d, y^{\prime}=d^{\prime}$ since

$$
\begin{aligned}
& \left\langle h_{i}, d\right\rangle=0 \quad \text { for } i=1, \ldots, l \\
& \langle c, d\rangle=1 \\
& \langle d, d\rangle=0
\end{aligned}
$$

Thus it is true for all $x^{\prime}, y^{\prime} \in L\left(A^{\prime}\right)$.
We now wish to obtain realisations of the remaining twisted Kac-Moody algebras of type $\tilde{C}_{l}^{\prime}$ for $l \geq 2$ and $\tilde{A}_{1}^{\prime}$. These will be obtained as fixed point subalgebras of the untwisted Kac-Moody algebra of type $\tilde{A}_{2 l}$ under the twisted graph automorphism $\tau$. We begin by recalling from Theorem 9.19 that the fixed point subalgebra of the finite dimensional Lie algebra $L\left(A_{2 l}\right)$ under its graph automorphism $\sigma$ is given by

$$
L\left(A_{2 l}\right)^{\sigma} \cong L\left(B_{l}\right)
$$



In order to show that $L\left(\tilde{A}_{2 l}\right)^{\tau} \cong L\left(\tilde{C}_{l}^{\prime}\right)$ if $l \geq 2$ and $L\left(\tilde{A}_{2}\right)^{\tau} \cong L\left(\tilde{A}_{1}^{\prime}\right)$ we shall compare the diagrams


$\tilde{A}_{2}$

$\tilde{A}_{1}^{\prime}$

We note that the numbering of the vertices of the graph for $\tilde{C}_{l}^{\prime}$ is not the same as the numbering previously used when $\tilde{C}_{l}^{\prime}$ was constructed from $C_{l}$ by adding an extra vertex labelled 0 . Here we are starting from the finite dimensional simple Lie algebra $B_{l}$ rather than $C_{l}$. In Theorem 17.17 (d) and (e) we obtained the roots of $\tilde{C}_{l}^{\prime}$ and $\tilde{A}_{1}^{\prime}$ in terms of those of $C_{l}$. For our present purpose we require these roots in terms of those of $B_{l}$.

We consider the diagram of $\tilde{C}_{l}^{\prime}$ labelled as follows.

and let $B_{l}$ be the subdiagram obtained by omitting vertex 0 and $C_{l}$ be the subdiagram obtained by omitting vertex $l$. We recall that

$$
\delta=\alpha_{0}+2 \alpha_{1}+\cdots+2 \alpha_{l-1}+2 \alpha_{l}
$$

The following lemma will be useful by relating the roots of $B_{l}$ and $C_{l}$.

Lemma 18.11 (i) Each long positive root of $C_{l}$ involves $\alpha_{0}$. Each short positive root of $B_{l}$ involves $\alpha_{l}$. There is a bijective correspondence $\alpha \leftrightarrow \beta$ between long positive roots of $C_{l}$ and short positive roots of $B_{l}$ satisfying $\alpha+2 \beta=\delta$.
(ii) There is a bijective correspondence $\alpha \leftrightarrow \beta$ between short positive roots of $C_{l}$ and long positive roots of $B_{l}$. If $\alpha, \beta$ do not involve $\alpha_{0}, \alpha_{l}$ respectively this correspondence is the identity map. If $\alpha$ involves $\alpha_{0}$ and $\beta$ involves $\alpha_{l}$ the correspondence is given by $\alpha+\beta=\delta$.

Proof. This follows immediately from expressing the roots of $B_{l}$ and $C_{l}$ in terms of the fundamental roots $\alpha_{1}, \ldots, \alpha_{l}$ and $\alpha_{0}, \alpha_{1}, \ldots, \alpha_{l-1}$ respectively.

Example The above bijection between $\Phi^{+}\left(C_{3}\right)$ and $\Phi^{+}\left(B_{3}\right)$ is as given below.


$$
\begin{array}{cc}
\frac{\Phi^{+}\left(C_{3}\right)}{\alpha_{0}} & \alpha_{1}+\Phi_{2}+\alpha_{3} \\
\alpha_{0}+2 \alpha_{1} & \alpha_{2}+\alpha_{3}
\end{array}
$$



By using Lemma 18.11 together with Theorem 17.17 we may express the real roots of $\tilde{C}_{l}^{\prime}$ in terms of the roots of $B_{l}$.

Proposition 18.12 (i) The real roots of $L\left(\tilde{C}_{l}^{\prime}\right), l \geq 2$, are

$$
\begin{aligned}
& \Phi_{\mathrm{Re}, \mathrm{~s}}=\left\{\alpha+r \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, \mathrm{i}}=\left\{\alpha+r \delta ; \alpha \in \Phi_{1}^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, \mathrm{l}}=\left\{2 \alpha+(2 r+1) \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}, r \in \mathbb{Z}\right\}
\end{aligned}
$$

where $\Phi_{\mathrm{Re}, \mathrm{s}}, \Phi_{\mathrm{Re}, \mathrm{i}}, \Phi_{\mathrm{Re}, \mathrm{l}}$ are the short, intermediate and long roots respectively, and $\Phi_{\mathrm{s}}^{0}, \Phi_{1}^{0}$ are the short and long roots of $B_{l}$.
(ii) The real roots of $L\left(\tilde{A}_{1}^{\prime}\right)$ are

$$
\begin{aligned}
& \Phi_{\mathrm{Re}, \mathrm{~s}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, 1}=\left\{2 \alpha+(2 r+1) \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\}
\end{aligned}
$$

where $\Phi^{0}$ is the root system of type $A_{1}$ obtained from the short fundamental root of $\tilde{A}_{1}^{\prime}$.

Proof. (i) We know from Theorem 17.17 that

$$
\begin{aligned}
& \Phi_{\mathrm{Re}, \mathrm{~s}}=\left\{\frac{1}{2}(\alpha+(2 r-1) \delta) ; \alpha \in \Phi_{l}^{0}\left(C_{l}\right), r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, \mathrm{i}}=\left\{\alpha+r \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}\left(C_{l}\right), r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, \mathrm{l}}=\left\{\alpha+2 r \delta ; \alpha \in \Phi_{1}^{0}\left(C_{l}\right), r \in \mathbb{Z}\right\}
\end{aligned}
$$

We make use of the bijection $\alpha \leftrightarrow \beta,-\alpha \leftrightarrow-\beta$ of Lemma 18.11 where $\alpha \in \Phi^{+}\left(C_{l}\right), \beta \in \Phi^{+}\left(B_{l}\right)$. For each $\alpha \in \Phi^{0}\left(C_{l}\right)$ we choose the corresponding
$\beta \in \Phi^{0}\left(B_{l}\right)$. First suppose $\alpha \in \Phi_{l}^{0}\left(C_{l}\right)$. The corresponding $\beta \in \Phi_{\mathrm{s}}^{0}\left(B_{l}\right)$ is given by $\alpha+2 \beta=\delta$ if $\alpha$ is positive and $\alpha+2 \beta=-\delta$ if $\alpha$ is negative. Thus

$$
\begin{aligned}
& \frac{1}{2}(\alpha+(2 r-1) \delta)= \begin{cases}-\beta+r \delta & \text { if } \alpha \text { is positive } \\
-\beta+(r-1) \delta & \text { if } \alpha \text { is negative }\end{cases} \\
& \alpha+2 r \delta= \begin{cases}-2 \beta+(2 r+1) \delta & \text { if } \alpha \text { is positive } \\
-2 \beta+(2 r-1) \delta & \text { if } \alpha \text { is negative }\end{cases}
\end{aligned}
$$

This gives the required formulae for $\Phi_{\mathrm{Re}, \mathrm{s}}$ and $\Phi_{\mathrm{Re}, \mathrm{l}}$ as $r$ runs through $\mathbb{Z}$. Next consider $\Phi_{\mathrm{Re}, \mathrm{i}}$. If $\alpha \in \Phi_{\mathrm{s}}^{0}\left(C_{l}\right)$ does not involve $\alpha_{0}$ we have $\beta=\alpha \in \Phi_{1}^{0}\left(B_{l}\right)$. If $\alpha \in \Phi_{\mathrm{s}}^{0}\left(C_{l}\right)$ does involve $\alpha_{0}$ the corresponding $\beta \in \Phi_{1}^{0}\left(B_{l}\right)$ satisfies $\alpha+\beta=\delta$ if $\alpha$ is positive and $\alpha+\beta=-\delta$ if $\alpha$ is negative. Then

$$
\begin{aligned}
& \alpha+r \delta=\beta+r \delta \quad \text { in the first case } \\
& \alpha+r \delta= \begin{cases}-\beta+(r+1) \delta & \text { if } \alpha \text { is positive } \\
-\beta+(r-1) \delta & \text { if } \alpha \text { is negative }\end{cases}
\end{aligned}
$$

in the second case. This gives the required formula for $\Phi_{\mathrm{Re}, \mathrm{i}}$ as $r$ runs through $\mathbb{Z}$.
(ii) In type $\tilde{A}_{1}^{\prime}$ the argument is exactly the same except that $\Phi^{0}=\Phi_{\mathrm{s}}^{0}$ has type $A_{1}$ and $\Phi_{1}^{0}$ is empty. Thus $\Phi_{\mathrm{Re}, \mathrm{i}}$ is empty in this case.

We shall next prove the analogue of Proposition 18.8 in our present case.
Proposition 18.13 Let $L^{0}$ be the simple Lie algebra of type $A_{2 l}$ and $\sigma$ be its graph automorphism of order 2. Let $\left(L^{0}\right)_{-1}$ be the eigenspace of $\sigma$ on $L^{0}$ with eigenvalue -1 . Then $L^{0}=\left(L^{0}\right)^{\sigma} \oplus\left(L^{0}\right)_{-1}^{1}$. The eigenspace $\left(L^{0}\right)_{-1}$ is an irreducible $\left(L^{0}\right)^{\sigma}$-module. The algebra $\left(L^{0}\right)^{\sigma}$ is isomorphic to $L\left(B_{l}\right)$ and $\left(L^{0}\right)_{-1}$ is isomorphic to its irreducible module $L\left(2 \omega_{1}\right)$.

Proof. It is clear that $\left(L^{0}\right)_{-1}$ is an $\left(L^{0}\right)^{\sigma}$-module and that $L^{0}=\left(L^{0}\right)^{\sigma} \oplus\left(L^{0}\right)_{-1}$. We have

$$
\begin{aligned}
& \operatorname{dim} L^{0}=\operatorname{dim} L\left(A_{2 l}\right)=2 l(2 l+2) \\
& \operatorname{dim}\left(L^{0}\right)^{\sigma}=\operatorname{dim} L\left(B_{l}\right)=l(2 l+1)
\end{aligned}
$$

Thus $\operatorname{dim}\left(L^{0}\right)_{-1}=\operatorname{dim} L^{0}-\operatorname{dim}\left(L^{0}\right)^{\sigma}=l(2 l+3)$. Let $\theta=\alpha_{1}+\cdots+\alpha_{2 l}$ be the highest root of $A_{2 l}$. Then a corresponding root vector $E_{\theta}$ lies in $\left(L^{0}\right)_{-1}$ and gives the highest weight of $\left(L^{0}\right)_{-1}$. It gives rise to the weight $2 \alpha_{1}+\cdots$ $+2 \alpha_{l}$ of $L\left(B_{l}\right)$. By considering the Cartan matrix of $B_{l}$ we see that $\alpha_{1}+\cdots$ $+\alpha_{l}=\omega_{1}$ and so the highest weight of the $L\left(B_{l}\right)$-module $\left(L^{0}\right)_{-1}$ is $2 \omega_{1}$.

Weyl's dimension formula shows that $\operatorname{dim} L\left(2 \omega_{1}\right)=l(2 l+3)$. Since this is also the dimension of $\left(L^{0}\right)_{-1}$ we see that $\left(L^{0}\right)_{-1}$ is an irreducible $L\left(B_{l}\right)$ module isomorphic to $L\left(2 \omega_{1}\right)$.

By using this result we can prove the analogue of Theorem 18.9 in the $A_{2 l}$ case.

Theorem 18.14 Let $L^{0}$ be a simple Lie algebra of type $A_{2 l}$ with $l \geq 2$ and $\sigma$ be its graph automorphism of order 2 . Let $\tau$ be the corresponding twisted graph automorphism of $\hat{\mathfrak{R}}\left(L^{0}\right)$. Then the fixed point subalgebra $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$ is isomorphic to $L\left(\tilde{C}_{l}^{\prime}\right)$.

When $l=1$ the fixed point subalgebra is isomorphic to $L\left(\tilde{A}_{1}^{\prime}\right)$.
Proof. The general idea of the proof is like that of Theorems 18.5 and 18.9. We aim to obtain the result by applying Proposition 14.15.

We first suppose $l \geq 2$. We have

$$
\begin{aligned}
& \hat{\mathfrak{R}}\left(L^{0}\right)=\sum_{k \in \mathbb{Z}}\left(t^{k} \otimes L^{0}\right) \oplus \mathbb{C} c \oplus \mathbb{C} d \\
& \hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}=\sum_{k \in \mathbb{Z}}\left(t^{2 k} \otimes\left(L^{0}\right)^{\sigma}\right)+\sum_{k \in \mathbb{Z}}\left(t^{2 k+1} \otimes\left(L^{0}\right)_{-1}\right) \oplus \mathbb{C} c \oplus \mathbb{C} d .
\end{aligned}
$$

Let $E_{1}, \ldots, E_{2 l}, F_{1}, \ldots, F_{2 l}, H_{1}, \ldots, H_{2 l}$ be standard generators of $L^{0}$. We wish to define analogous elements $e_{0}, e_{1}, \ldots, e_{l}, f_{0}, f_{1}, \ldots, f_{l}, h_{0}, h_{1}, \ldots, h_{l}$ in $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$. These are chosen as follows.

$$
\begin{array}{llll}
e_{1}=1 \otimes\left(E_{1}+E_{2 l}\right), & \ldots, & e_{l-1}=1 \otimes\left(E_{l-1}+E_{l+2}\right), & e_{l}=1 \otimes \sqrt{ } 2\left(E_{l}+E_{l+1}\right) \\
f_{1}=1 \otimes\left(F_{1}+F_{2 l}\right), & \ldots, & f_{l-1}=1 \otimes\left(F_{l-1}+F_{l+2}\right), & f_{l}=1 \otimes \sqrt{ } 2\left(F_{l}+F_{l+1}\right) \\
h_{1}=1 \otimes\left(H_{1}+H_{2 l}\right), & \ldots, & h_{l-1}=1 \otimes\left(H_{l-1}+H_{l+2}\right), & h_{l}=1 \otimes 2\left(H_{l}+H_{l+1}\right) \\
e_{0}=t \otimes F_{\theta}, & f_{0}=t^{-1} \otimes E_{\theta}, & h_{0}=-\left(1 \otimes H_{\theta}\right)+c
\end{array}
$$

where $\theta$ is the highest root of $L^{0}$.
Let $H \subset \hat{\mathfrak{Z}}\left(L^{0}\right)$ be given by $H=\left(1 \otimes H^{0}\right) \oplus \mathbb{C} c \oplus \mathbb{C} d$. We define maps $\alpha_{0}, \alpha_{1}, \ldots, \alpha_{l}: H^{\sigma} \rightarrow \mathbb{C}$. For $i=1, \ldots, l$ these are the restrictions of the roots $\alpha_{i}: H \rightarrow \mathbb{C}$. For $i=0$ we define $\alpha_{0}=-\theta+\delta$.

Let $\Pi^{\mathrm{v}}=\left\{h_{0}, h_{1}, \ldots, h_{l}\right\} \subset H^{\sigma}$ and $\Pi=\left\{\alpha_{0}, \alpha_{1}, \ldots, \alpha_{l}\right\} \subset\left(H^{\sigma}\right)^{*}$. We show that $\left(H^{\sigma}, \Pi, \Pi^{v}\right)$ is a realisation of the Cartan matrix $A^{\prime}$ of type $\tilde{C}_{l}^{\prime}$.

We know from Theorem 9.19 that $\alpha_{j}\left(h_{i}\right)=A_{i j}^{\prime}$ for $i, j \in\{1, \ldots, l\}$. This $l \times l$ matrix is the Cartan matrix of type $B_{l}$. In particular it is non-singular. Thus $h_{1}, \ldots, h_{l}$ and $\alpha_{1}, \ldots, \alpha_{l}$ are linearly independent. Now $h_{0}$ involves $c$ whereas $h_{1}, \ldots, h_{l}$ do not, thus $h_{0}, h_{1}, \ldots, h_{l}$ are linearly independent.

Also we have $\alpha_{0}(d) \neq 0$ and $\alpha_{i}(d)=0$ for $i=1, \ldots, l$ thus $\alpha_{0}, \alpha_{1}, \ldots, \alpha_{l}$ are linearly independent. We must show in addition

$$
\begin{array}{ll}
\alpha_{j}\left(h_{0}\right)=A_{0 j}^{\prime} & j=1, \ldots, l \\
\alpha_{0}\left(h_{i}\right)=A_{i 0}^{\prime} & i=1, \ldots, l \\
\alpha_{0}\left(h_{0}\right)=2 . &
\end{array}
$$

We recall that the integers $a_{i}, c_{i}$ associated with the affine Cartan matrix $A^{\prime}$ are


In particular we have $a_{0}=1, c_{0}=2$. (The change of labelling explains the fact that $c_{0}$ is not 1 , as it usually is.) We note that

$$
\begin{aligned}
\sum_{i=0}^{l} a_{i} \alpha_{i} & =\delta \\
\sum_{i=0}^{l} c_{i} h_{i} & =2 c
\end{aligned}
$$

We then have

$$
\begin{aligned}
& \alpha_{0}\left(h_{i}\right)=-\sum_{j=1}^{l} a_{j} \alpha_{j}\left(h_{i}\right)=-\sum_{j=1}^{l} A_{i j}^{\prime} a_{j}=A_{i 0}^{\prime} a_{0}=A_{i 0}^{\prime} \\
& \alpha_{j}\left(h_{0}\right)=-\frac{1}{2} \sum_{i=1}^{l} c_{i} \alpha_{j}\left(h_{i}\right)=-\frac{1}{2} \sum_{i=1}^{l} c_{i} A_{i j}^{\prime}=\frac{1}{2} c_{0} A_{0 j}^{\prime}=A_{0 j}^{\prime} \\
& \alpha_{0}\left(h_{0}\right)=(-\theta+\delta)\left(\left(1 \otimes-H_{\theta}\right)+c\right)=\theta\left(H_{\theta}\right)=2
\end{aligned}
$$

We observe that $A^{\prime}$ is an $(l+1) \times(l+1)$ matrix of rank $l$ and that $\operatorname{dim} H^{\sigma}=$ $l+2$. Thus we have shown that $\left(H^{\sigma}, \Pi, \Pi^{v}\right)$ is a realisation of $A^{\prime}$.

We next verify the relations

$$
\left[h_{i} e_{j}\right]=A_{i j}^{\prime} e_{j} \quad\left[h_{i} f_{j}\right]=-A_{i j}^{\prime} f_{j}
$$

for $i, j \in\{0,1, \ldots, l\}$. We know this already for $i, j \in\{1, \ldots, l\}$ by Theorem 9.19. Thus we must verify

$$
\begin{array}{ll}
{\left[h_{0} e_{j}\right]=A_{0 j}^{\prime} e_{j}} & {\left[h_{0} f_{j}\right]=-A_{0 j}^{\prime} f_{j}} \\
{\left[h_{i} e_{0}\right]=A_{i 0}^{\prime} e_{0}} & {\left[h_{i} f_{0}\right]=-A_{i 0}^{\prime} f_{0}} \\
{\left[h_{0} e_{0}\right]=2 e_{0}} & {\left[h_{0} f_{0}\right]=-2 f_{0}}
\end{array}
$$

Now

$$
\left[h_{0} e_{j}\right]=-\frac{1}{2} \sum_{i=1}^{l} c_{i}\left[h_{i} e_{j}\right]=\left(-\frac{1}{2} \sum_{i=1}^{l} c_{i} A_{i j}^{\prime}\right) e_{j}=A_{0 j}^{\prime} e_{j}
$$

and similarly $\left[h_{0} f_{j}\right]=-A_{0 j}^{\prime} f_{j}$. Also

$$
\begin{aligned}
{\left[h_{i} e_{0}\right] } & =\left[h_{i}, t \otimes F_{\theta}\right]=t \otimes\left(-\theta\left(h_{i}\right) F_{\theta}\right)=-\theta\left(h_{i}\right) e_{0} \\
& =\left(-\sum_{j=1}^{l} a_{j} \alpha_{j}\left(h_{i}\right)\right) e_{0}=\left(-\sum_{j=1}^{l} A_{i j}^{\prime} a_{j}\right) e_{0}=A_{i 0}^{\prime} e_{0}
\end{aligned}
$$

Similarly we have $\left[h_{i} f_{0}\right]=-A_{i 0}^{\prime} f_{0}$. We also have

$$
\begin{aligned}
{\left[h_{0} e_{0}\right] } & =\left[-\left(1 \otimes H_{\theta}\right)+c, t \otimes F_{\theta}\right]=-t \otimes\left[H_{\theta} F_{\theta}\right] \\
& =2 t \otimes F_{\theta}=2 e_{0}
\end{aligned}
$$

and similarly $\left[h_{0} f_{0}\right]=-2 f_{0}$. Finally we have

$$
\begin{aligned}
& {\left[c e_{i}\right]=\alpha_{i}(c) e_{i}=0 \quad i=0,1, \ldots, l} \\
& {\left[c f_{i}\right]=-\alpha_{i}(c) f_{i}=0 \quad i=0,1, \ldots, l} \\
& {\left[d e_{i}\right]=\alpha_{i}(d) e_{i}=0 \quad i=1, \ldots, l} \\
& {\left[d f_{i}\right]=-\alpha_{i}(d) f_{i}=0 \quad i=1, \ldots, l} \\
& {\left[d e_{0}\right]=\alpha_{0}(d) e_{0}=e_{0}} \\
& {\left[d f_{0}\right]=-\alpha_{0}(d) f_{0}=-f_{0} .}
\end{aligned}
$$

Since $H^{\tau}=\left(1 \otimes\left(H^{0}\right)^{\sigma}\right) \oplus \mathbb{C} c \oplus \mathbb{C} d$ we have verified all relations necessary for the application of Proposition 14.15.

We next show that the elements $e_{0}, e_{1}, \ldots, e_{l}, f_{0}, f_{1}, \ldots, f_{l}$ together with $H^{\tau}$ generate $\hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}$. By Theorem $9.19 e_{1}, \ldots, e_{l}, f_{1}, \ldots, f_{l}$ generate $\left(L^{0}\right)^{\sigma}$. Since $\hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}=\sum_{k \in \mathbb{Z}}\left(t^{2 k} \otimes\left(L^{0}\right)^{\sigma}\right) \oplus \sum_{k \in \mathbb{Z}}\left(t^{2 k+1} \otimes\left(L^{0}\right)_{-1}\right) \oplus \mathbb{C} c \oplus \mathbb{C} d$ it is sufficient to show that the subspaces

$$
\left(t^{2 k} \otimes\left(L^{0}\right)^{\sigma}\right) \quad \text { for } k \neq 0 \quad \text { and }\left(t^{2 k+1} \otimes\left(L^{0}\right)_{-1}\right)
$$

lie in the subalgebra $M$ generated by $e_{0}, e_{1}, \ldots, e_{l}, f_{0}, f_{1}, \ldots, f_{l}, H^{\tau}$.
Now $e_{0}=t \otimes F_{\theta}$ lies in $M$ and $F_{\theta} \in\left(L^{0}\right)_{-1}$. If $x \in\left(L^{0}\right)^{\sigma}, y \in\left(L^{0}\right)_{-1}$ then

$$
[1 \otimes x, t \otimes y]=t \otimes[x y] \in t \otimes\left(L^{0}\right)_{-1}
$$

Thus the elements $y \in\left(L^{0}\right)_{-1}$ for which $t \otimes y \in M$ form an $\left(L^{0}\right)^{\sigma}$-submodule of $\left(L^{0}\right)_{-1}$. This submodule contains $F_{\theta}$ so is non-zero. Since $\left(L^{0}\right)_{-1}$ is an
irreducible $\left(L^{0}\right)^{\sigma}$-module by Proposition 18.13 this submodule is the whole of $\left(L^{0}\right)_{-1}$. Thus $t \otimes\left(L^{0}\right)_{-1}$ lies in $M$.

Now we can find elements $x, y \in\left(L^{0}\right)_{-1}$ with $[x y] \neq 0$. (For example, $x=$ $F_{\alpha_{l}}-F_{\alpha_{l+1}}, y=E_{\alpha_{l}+\alpha_{l+1}}$. .) Thus $[t \otimes x, t \otimes y]=t^{2} \otimes[x y]$ is a non-zero element of $M$. However, the set of $z \in\left(L^{0}\right)^{\sigma}$ for which $t^{2} \otimes z \in M$ is an ideal of $\left(L^{0}\right)^{\sigma}$ since

$$
\left[t^{2} \otimes z, 1 \otimes w\right]=t^{2} \otimes[z w] \quad \text { for } w \in\left(L^{0}\right)^{\sigma}
$$

Since $[x y] \in\left(L^{0}\right)^{\sigma}$ this is a non-zero ideal of $\left(L^{0}\right)^{\sigma}$ and since $\left(L^{0}\right)^{\sigma}$ is a simple Lie algebra it is the whole of $\left(L^{0}\right)^{\sigma}$. Thus $t^{2} \otimes\left(L^{0}\right)^{\sigma}$ lies in $M$.

Now the relations

$$
\begin{aligned}
& {\left[t^{2} \otimes x, t^{2 k} \otimes y\right]=t^{2 k+2} \otimes[x y], \quad x, y \in\left(L^{0}\right)^{\sigma}} \\
& {\left[t^{2} \otimes x, t^{2 k+1} \otimes y\right]=t^{2 k+3} \otimes[x y], \quad x \in\left(L^{0}\right)^{\sigma}, y \in\left(L^{0}\right)_{-1}}
\end{aligned}
$$

can be used to show by induction on $k$ that $t^{2 k} \otimes\left(L^{0}\right)^{\sigma} \subset M$ and

$$
t^{2 k+1} \otimes\left(L^{0}\right)_{-1} \subset M \quad \text { when } k>0
$$

Starting with $f_{0}$ instead of $e_{0}$ will similarly show this for $k<0$. Thus $M=\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$ as required.

Finally we must show that $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$ has no non-zero ideal $J$ with $J \cap H^{\tau}=O$. To see this we decompose $\hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}$ into weight spaces with respect to $H^{\tau}$. Any non-zero ideal $J$ with $J \cap H^{\tau}=O$ must have non-zero intersection with one of these weight spaces, by Lemma 14.12. Let $x$ be a non-zero element in such an intersection. Then there exists $y$ in the weight space corresponding to the negative of this weight such that $[x y] \neq 0$, by Corollary 16.5. But then $[x y] \in J \cap H^{\tau}$ and so $J \cap H^{\tau} \neq O$, a contradiction. Hence $J=O$.

We have now verified all the hypotheses of the recognition theorem Proposition 14.15 and so can conclude that $\hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}$ is isomorphic to $L\left(\tilde{C}_{l}^{\prime}\right)$.

We now consider the case $l=1$. This time the graphs are

and the Cartan matrix $A^{\prime}$ of $\tilde{A}_{1}^{\prime}$ is

$$
A^{\prime}=\frac{\left(\begin{array}{rr}
2 & -1 \\
-4 & 2
\end{array}\right)}{0} \begin{aligned}
& 0 \\
& 0
\end{aligned}
$$

The elements $e_{0}, e_{1}, f_{0}, f_{1}, h_{0}, h_{1}$ of $\hat{\mathbb{R}}\left(L^{0}\right)^{\tau}$ are

$$
\begin{aligned}
& e_{1}=1 \otimes \sqrt{ } 2\left(E_{1}+E_{2}\right), \quad f_{1}=1 \otimes \sqrt{ } 2\left(F_{1}+F_{2}\right), \quad h_{1}=1 \otimes 2\left(H_{1}+H_{2}\right) \\
& e_{0}=t \otimes F_{\theta}, \quad f_{0}=t^{-1} \otimes E_{\theta}, \quad h_{0}=-\left(1 \otimes H_{\theta}\right)+c
\end{aligned}
$$

where $\theta=\alpha_{1}+\alpha_{2}$ is the highest root of $A_{2}$. We have

$$
\stackrel{\mathrm{v}}{\Pi}=\left\{h_{0}, h_{1}\right\} \quad \Pi=\left\{\alpha_{0}, \alpha_{1}\right\}
$$

where $\alpha_{0}=-\theta+\delta$ and $\alpha_{1} \in\left(H^{\sigma}\right)^{*}$ is the restriction of $\alpha_{1} \in H^{*}$. The integers $a_{0}, a_{1}, c_{0}, c_{1}$ for $\tilde{A}_{1}^{\prime}$ are

$$
a_{0}=1, \quad a_{1}=2, \quad c_{0}=2, \quad c_{1}=1
$$

We have

$$
\begin{aligned}
& a_{0} \alpha_{0}+a_{1} \alpha_{1}=\delta \\
& c_{0} h_{0}+c_{1} h_{1}=2 c
\end{aligned}
$$

We can then check that $\left(H^{\sigma}, \Pi, \Pi^{v}\right)$ is a realisation of $A^{\prime}$. We also check the relations

$$
\begin{array}{llll}
{\left[h_{0} e_{0}\right]=2 e_{0},} & {\left[h_{0} e_{1}\right]=-e_{1},} & {\left[h_{1} e_{0}\right]=-4 e_{0},} & {\left[h_{1} e_{1}\right]=2 e_{1}} \\
{\left[h_{0} f_{0}\right]=-2 f_{0},} & {\left[h_{0} f_{1}\right]=f_{1},} & {\left[h_{1} f_{0}\right]=4 f_{0},} & {\left[h_{1} f_{1}\right]=-2 f_{1}}
\end{array}
$$

The facts that $H^{\tau}, e_{0}, e_{1}, f_{0}, f_{1}$ generate $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$ and that $\hat{\mathfrak{R}}\left(L^{0}\right)^{\tau}$ has no non-zero ideal $J$ with $J \cap H^{\tau}=O$ are proved just as before. Thus applying Proposition 14.15 shows that $\hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}$ is isomorphic to $L\left(\tilde{A}_{1}^{\prime}\right)$.

We shall describe explicitly the weight space decomposition of $\hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}$ with respect to $H^{\tau}$. We recall from Proposition 9.18 that there is a bijective correspondence between roots of $\left(L^{0}\right)^{\sigma}=L\left(B_{l}\right)$ and equivalence classes of roots of $L^{0}=L\left(A_{2 l}\right)$. Each equivalence class has 2 or 3 elements. Equivalence classes with 2 elements have form $(\alpha, \beta)$ where $\sigma(\alpha)=\beta, \sigma(\beta)=\alpha$ and $\alpha+\beta$ is not a root. Equivalence classes with 3 elements have form $(\alpha, \beta, \alpha+\beta)$ where $\sigma(\alpha)=\beta, \sigma(\beta)=\alpha$ and $\sigma(\alpha+\beta)=\alpha+\beta$. Equivalence classes with 2 elements correspond to long roots of $B_{l}$ and equivalence classes with 3 elements correspond to short roots of $B_{l}$. For each 2-element equivalence class we can choose root vectors $E_{\alpha}, E_{\beta}$ with $\sigma\left(E_{\alpha}\right)=E_{\beta}$. For each 3-element equivalence class we choose root vectors $E_{\alpha}, E_{\beta}, E_{\alpha+\beta}$ with $\sigma\left(E_{\alpha}\right)=E_{\beta}$ and $\left[E_{\alpha} E_{\beta}\right]=E_{\alpha+\beta}$. Then

$$
\sigma\left(E_{\alpha+\beta}\right)=\sigma\left[E_{\alpha} E_{\beta}\right]=\left[E_{\beta} E_{\alpha}\right]=-E_{\alpha+\beta}
$$

thus $E_{\alpha+\beta} \in\left(L^{0}\right)_{-1}$.

The Lie algebra $\hat{\Omega}\left(L^{0}\right)^{\tau}$ is the direct sum of $H^{\tau}$ and the following weight spaces.
$t^{2 k} \otimes\left(H^{0}\right)^{\sigma} \quad$ with weight $2 k \delta$
$t^{2 k+1} \otimes\left(H^{0}\right)_{-1} \quad$ with weight $(2 k+1) \delta$
$t^{2 k} \otimes \mathbb{C}\left(E_{\alpha}+E_{\beta}\right) \quad$ with weight $\alpha+2 k \delta$ for each 2 element equivalence class $(\alpha, \beta)$

$$
\begin{array}{ll}
t^{2 k+1} \otimes \mathbb{C}\left(E_{\alpha}-E_{\beta}\right) & \begin{array}{l}
\text { with weight } \alpha+(2 k+1) \delta \text { for each } 2 \text { element } \\
\text { equivalence class }(\alpha, \beta)
\end{array} \\
t^{2 k} \otimes \mathbb{C}\left(E_{\alpha}+E_{\beta}\right) & \begin{array}{l}
\text { with weight } \alpha+2 k \delta \text { for each } 3 \text { element equivalence } \\
\text { class }(\alpha, \beta, \alpha+\beta)
\end{array}
\end{array}
$$

$t^{2 k+1} \otimes \mathbb{C}\left(E_{\alpha}-E_{\beta}\right) \quad$ with weight $\alpha+(2 k+1) \delta$ for each 3 element equivalence class $(\alpha, \beta, \alpha+\beta)$
$t^{2 k+1} \otimes \mathbb{C} E_{\alpha+\beta} \quad$ with weight $2 \alpha+(2 k+1) \delta$ for each 3 element equivalence $\operatorname{class}(\alpha, \beta, \alpha+\beta)$.

The weights listed above correspond to the roots of $L\left(\tilde{C}_{l}^{\prime}\right)$ as described in Proposition 18.12.

In the case $l=1$ the weight spaces of $L\left(\tilde{A}_{2}\right)^{\tau}$ are

$$
\begin{array}{ll}
t^{2 k} \otimes \mathbb{C}\left(H_{1}+H_{2}\right) \quad \text { with weight } 2 k \delta \\
t^{2 k+1} \otimes \mathbb{C}\left(H_{1}-H_{2}\right) \quad \text { with weight }(2 k+1) \delta \\
t^{2 k} \otimes \mathbb{C}\left(E_{1}+E_{2}\right) \quad \text { with weight } \alpha_{1}+2 k \delta \\
t^{2 k+1} \otimes \mathbb{C}\left(E_{1}-E_{2}\right) \quad \text { with weight } \alpha_{1}+(2 k+1) \delta \\
t^{2 k+1} \otimes \mathbb{C} E_{\alpha_{1}+\alpha_{2}} \quad \text { with weight } 2 \alpha_{1}+(2 k+1) \delta \\
t^{2 k} \otimes \mathbb{C}\left(F_{1}+F_{2}\right) \quad \text { with weight }-\alpha_{1}+2 k \delta \\
t^{2 k+1} \otimes \mathbb{C}\left(F_{1}-F_{2}\right) \quad \text { with weight }-\alpha_{1}+(2 k+1) \delta \\
t^{2 k+1} \otimes \mathbb{C} F_{\alpha_{1}+\alpha_{2}} \quad \text { with weight-2 } 2+(2 k+1) \delta .
\end{array}
$$

Corollary 18.15 (i) The multiplicities of the imaginary roots $k \delta$ of $\tilde{C}_{l}^{\prime}$ are equal to $l$.
(ii) The multiplicities of the imaginary roots $k \delta$ of $\tilde{A}_{1}^{\prime}$ are equal to 1 .

Proof. (i) The multiplicity of $2 k \delta$ is $\operatorname{dim}\left(H^{0}\right)^{\sigma}$, which is equal to $l$. The multiplicity of $(2 k+1) \delta$ is $\operatorname{dim}\left(H^{0}\right)_{-1}$, which is also equal to $l$.
(ii) The same applies to $\tilde{A}_{1}^{\prime}$ when $l=1$.

We note that the isomorphism between $L\left(A^{\prime}\right)$ and $\hat{\mathfrak{R}}\left(A_{2 l}\right)^{\tau}$ does not map the standard invariant form $\langle,\rangle^{\prime}$ on $L\left(A^{\prime}\right)$ to the restriction of the standard invariant form $\langle$,$\rangle on \hat{\mathfrak{Z}}\left(A_{2 l}\right)$. For $h_{l} \in L\left(A^{\prime}\right)$ corresponds to $1 \otimes 2\left(H_{l}+H_{l+1}\right)$ in $\hat{\mathfrak{R}}\left(A_{2 l}\right)$. We have

$$
\begin{aligned}
& \left\langle h_{l}, h_{l}\right\rangle^{\prime}=2 \frac{a_{l}}{c_{l}}=4 \\
& \left\langle 1 \otimes 2\left(H_{l}+H_{l+1}\right), 1 \otimes 2\left(H_{l}+H_{l+1}\right)\right\rangle=4\left\langle H_{l}+H_{l+1}, H_{l}+H_{l+1}\right\rangle=8 .
\end{aligned}
$$

Thus the form is not preserved by the isomorphism.
Also the canonical central element $c \in \hat{\mathcal{R}}\left(A_{2 l}\right)^{\tau}$ does not map to the canonical central element $c^{\prime} \in L\left(A^{\prime}\right)$. Since we showed that $\sum_{i=0}^{l} c_{i} h_{i}=2 c$ it follows that $2 c$ corresponds to $c^{\prime}$ under our isomorphism.

## Comments on notation

An alternative notation is sometimes given to the affine Kac-Moody algebras of twisted type, based on the results of this chapter. The twisted affine algebra can be specified by the type of the untwisted affine algebra from which it is obtained, together with the order of the automorphism of which it is the fixed point subalgebra. This is the notation used by Kac in his book Infinite Dimensional Lie Algebras. The alternative notation in each case is shown below.

| $\tilde{B}_{l}^{\mathrm{t}}$ | ${ }^{2} \tilde{A}_{2 l-1}$ | $l \geq 3$ |
| :--- | :--- | :--- |
| $\tilde{C}_{l}^{\mathrm{t}}$ | ${ }^{2} \tilde{D}_{l+1}$ | $l \geq 2$ |
| $\tilde{F}_{4}^{\mathrm{t}}$ | ${ }^{2} \tilde{E}_{6}$ |  |
| $\tilde{G}_{2}^{\mathrm{t}}$ | ${ }^{3} \tilde{D}_{4}$ |  |
| $\tilde{C}_{l}^{\prime}$ | ${ }^{2} \tilde{A}_{2 l}$ | $l \geq 2$ |
| $\tilde{A}_{1}^{\prime}$ | ${ }^{2} \tilde{A}_{2}$ |  |

## 19

## Some representations of symmetrisable Kac-Moody algebras

### 19.1 The category $\mathcal{O}$ of $L(A)$-modules

We now turn to the representation theory of Kac-Moody algebras. We shall not consider arbitrary representations, but restrict attention to those in the category $\mathcal{O}$ introduced by Bernstein, Gelfand and Gelfand. Let

$$
L(A)=N^{-} \oplus H \oplus N
$$

be a Kac-Moody algebra and $V$ be an $L(A)$-module. We say that $V$ is an object in the category $\mathcal{O}$ if the following conditions are satisfied:
(i) $V=\bigoplus_{\lambda \in H^{*}} V_{\lambda}$ where $V_{\lambda}=\{v \in V ; x v=\lambda(x) v$ for all $x \in H\}$
(ii) $\operatorname{dim} V_{\lambda}$ is finite for each $\lambda \in H^{*}$
(iii) there exists a finite set $\lambda_{1}, \ldots, \lambda_{s} \in H^{*}$ such that each $\lambda$ with $V_{\lambda} \neq O$ satisfies $\lambda<\lambda_{i}$ for some $i \in\{1, \ldots, s\}$.

The morphisms in category $\mathcal{O}$ are the homomorphisms of $L(A)$-modules.
Thus each module in $\mathcal{O}$ is a direct sum of its weight spaces and these weight spaces are finite dimensional. Moreover all the weights are bounded above by finitely many elements of $H^{*}$.

We now give some examples of modules in category $\mathcal{O}$. For each $\lambda \in H^{*}$ we may define the Verma module $M(\lambda)$ with highest weight $\lambda$. This is defined in a manner analogous to that in which we defined Verma modules for finite dimensional Lie algebras in Section 10.1. Let $\mathfrak{l}(L(A))$ be the universal enveloping algebra of $L(A)$ and $K_{\lambda}$ be the left ideal of $\mathfrak{H}(L(A))$ generated by $N$ and $x-\lambda(x)$ for all $x \in H$. Thus

$$
K_{\lambda}=\mathfrak{U}(L(A)) N+\sum_{x \in H} \mathfrak{H}(L(A))(x-\lambda(x)) .
$$

Then $M(\lambda)=\mathfrak{U}(L(A)) / K_{\lambda}$ is an $L(A)$-module called the Verma module with highest weight $\lambda$. Let $m_{\lambda} \in M(\lambda)$ be defined by $m_{\lambda}=1+K_{\lambda}$. Then, just as in Theorem 10.6, we see that each element of $M(\lambda)$ is uniquely expressible in
the form $u m_{\lambda}$ for some $u \in \mathfrak{U}\left(N^{-}\right)$. Also, as in Theorem 10.7, we have

$$
\begin{aligned}
& M(\lambda)=\bigoplus_{\mu \in H^{*}} M(\lambda)_{\mu} \\
& M(\lambda)_{\mu} \neq O \quad \text { if and only if } \mu \prec \lambda \\
& \operatorname{dim} M(\lambda)_{\mu}=\mathfrak{B}(\lambda-\mu) .
\end{aligned}
$$

This shows that $M(\lambda) \in \mathcal{O}$. The finite set of weights giving an upper bound for all weights can be taken in this case to have just one element $\lambda$.

Lemma 19.1 (i) If $V \in \mathcal{O}$ and $U$ is a submodule of $V$ then $U \in \mathcal{O}$ and $V / U \in \mathcal{O}$. (ii) If $V_{1}, V_{2} \in \mathcal{O}$ then $V_{1} \oplus V_{2} \in \mathcal{O}$ and $V_{1} \otimes V_{2} \in \mathcal{O}$.

Proof. (i) We have $V=\bigoplus_{\lambda \in H^{*}} V_{\lambda}$. The argument of Theorem 10.9 shows that $U_{\lambda}=U \cap V_{\lambda}$ and $U=\bigoplus_{\lambda \in H^{*}} U_{\lambda}$. It follows that $U \in \mathcal{O}$. Moreover we have $(V / U)_{\lambda}=V_{\lambda} / U_{\lambda}$ and $V / U=\bigoplus_{\lambda \in H^{*}}(V / U)_{\lambda}$. It follows that $V / U \in \mathcal{O}$.
(ii) We have $\left(V_{1} \oplus V_{2}\right)_{\lambda}=\left(V_{1}\right)_{\lambda} \oplus\left(V_{2}\right)_{\lambda}$ and $V_{1} \oplus V_{2}=\bigoplus_{\lambda \in H^{*}}\left(V_{1} \oplus V_{2}\right)_{\lambda}$. It follows that $V_{1} \oplus V_{2} \in \mathcal{O}$.

Now consider $V_{1} \otimes V_{2}$. We have

$$
V_{1}=\bigoplus_{\lambda_{1} \in H^{*}}\left(V_{1}\right)_{\lambda_{1}}, \quad V_{2}=\bigoplus_{\lambda_{2} \in H^{*}}\left(V_{2}\right)_{\lambda_{2}}
$$

thus

$$
V_{1} \otimes V_{2}=\bigoplus_{\lambda_{1}, \lambda_{2}}\left(\left(V_{1}\right)_{\lambda_{1}} \otimes\left(V_{2}\right)_{\lambda_{2}}\right)
$$

Now $\left(V_{1}\right)_{\lambda_{1}} \otimes\left(V_{2}\right)_{\lambda_{2}} \subset\left(V_{1} \otimes V_{2}\right)_{\lambda_{1}+\lambda_{2}}$. Hence $V_{1} \otimes V_{2}=\bigoplus_{\lambda \in H^{*}}\left(V_{1} \otimes V_{2}\right)_{\lambda}$ where

$$
\left(V_{1} \otimes V_{2}\right)_{\lambda}=\sum_{\substack{\lambda_{1}, \lambda_{2} \\ \lambda_{1}+\lambda_{2}=\lambda}}\left(\left(V_{1}\right)_{\lambda_{1}} \otimes\left(V_{2}\right)_{\lambda_{2}}\right)
$$

Now there exist $\xi_{i} \in H^{*}, i=1, \ldots, s_{1}$, such that $\lambda_{1} \prec \xi_{i}$ for some $i$. Also there exist $\eta_{j} \in H^{*}, j=1, \ldots, s_{2}$ such that $\lambda_{2} \prec \eta_{j}$ for some $j$. Thus $\lambda=$ $\lambda_{1}+\lambda_{2} \prec \xi_{i}+\eta_{j}$ for some pair $(i, j)$. We have

$$
\left(\xi_{i}+\eta_{j}\right)-\lambda=\left(\xi_{i}-\lambda_{1}\right)+\left(\eta_{j}-\lambda_{2}\right)
$$

The expressions $\left(\xi_{i}+\eta_{j}\right)-\lambda, \xi_{i}-\lambda_{1}, \eta_{j}-\lambda_{2}$ are all non-negative integral combinations of the fundamental roots. Thus for given $i, j$ the $\left(\xi_{i}+\eta_{j}\right)-\lambda$ has only finitely many such decompositions. It follows that for each $\lambda$ with $\left(V_{1} \otimes V_{2}\right)_{\lambda} \neq O$ there exist only finitely many pairs $\lambda_{1}, \lambda_{2}$ with $\lambda_{1}+\lambda_{2}=\lambda,\left(V_{1}\right)_{\lambda_{1}} \neq O,\left(V_{2}\right)_{\lambda_{2}} \neq O$. It follows from this that $V_{1} \otimes V_{2} \in \mathcal{O}$.

Now each $L(A)$-module $V \in \mathcal{O}$ admits a character ch $V$. We recall from Section 12.1 that ch $V$ is the function from $H^{*}$ to $\mathbb{Z}$ defined by

$$
(\operatorname{ch} V)(\lambda)=\operatorname{dim} V_{\lambda} .
$$

We also recall from Section 12.1 the definition of the ring $\Re$ of functions from $H^{*}$ to $\mathbb{Z}$. A function $f: H^{*} \rightarrow \mathbb{Z}$ lies in $\Re$ if there exists a finite set $\lambda_{1}, \ldots, \lambda_{s} \in H^{*}$ such that

$$
\operatorname{Supp} f \subset S\left(\lambda_{1}\right) \cup \cdots \cup S\left(\lambda_{s}\right)
$$

where $S(\lambda)=\operatorname{Supp}(\operatorname{ch} M(\lambda))$. It follows from the definition of category $\mathcal{O}$ that ch $V \in \mathfrak{R}$ for all $V \in \mathcal{O}$.

In Proposition 12.4 we obtained a formula for the character of a Verma module for a finite dimensional semisimple Lie algebra. We now generalise this result to Verma modules for Kac-Moody algebras. We recall that the function $e_{\lambda}: H^{*} \rightarrow \mathbb{Z}$ was defined by $e_{\lambda}(\lambda)=1$ and $e_{\lambda}(\mu)=0$ if $\lambda \neq \mu$. The characteristic functions $e_{\lambda}$ lie in $\mathfrak{R}$ and any function $f \in \mathfrak{R}$ can be written in the form

$$
f=\sum_{\lambda \in H^{*}} f(\lambda) e_{\lambda}
$$

where the sum may be infinite.

Proposition 19.2 Let $M(\lambda)$ be a Verma module for the Kac-Moody algebra $L(A)$. Then

$$
\operatorname{ch} M(\lambda)=\frac{e_{\lambda}}{\prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)^{m_{\alpha}}}
$$

where $m_{\alpha}$ is the multiplicity of $\alpha$.
Proof. We use the fact that the map $u \rightarrow u m_{\lambda}$ is a bijection between $\mathfrak{H}\left(N^{-}\right)$ and $M(\lambda)$. This bijection maps the weight space $\mathfrak{l}\left(N^{-}\right)_{-\mu}$ to the weight space $M(\lambda)_{\lambda-\mu}$.

For each $\alpha \in \Phi^{+}$we have $\operatorname{dim}\left(N^{-}\right)_{-\alpha}=m_{\alpha}$. Let $e_{-\alpha, i}, 1 \leq i \leq m_{\alpha}$, be a basis of $\left(N^{-}\right)_{-\alpha}$. We choose an order on these basis elements for all $\alpha, i$. We then obtain a PBW-basis of $\mathfrak{H}\left(N^{-}\right)$consisting of all products

$$
\prod_{\alpha} \prod_{i=1}^{m_{\alpha}} e_{-\alpha, i}^{n_{\alpha, i}}
$$

with $n_{\alpha, i} \in \mathbb{Z}$ and $n_{\alpha, i} \geq 0$. Thus the weight space $\mathfrak{H}\left(N^{-}\right)_{-\mu}$ has a basis consisting of the above elements which satisfy

$$
\sum_{\alpha \in \Phi^{+}}\left(\sum_{i=1}^{m_{\alpha}} n_{\alpha, i}\right) \alpha=\mu
$$

This shows that the character of $\mathfrak{l}\left(N^{-}\right)$is

$$
\operatorname{ch} \mathfrak{U}\left(N^{-}\right)=\prod_{\alpha \in \Phi^{+}}\left(1+e_{-\alpha}+e_{-\alpha}^{2}+\cdots\right)^{m_{\alpha}}
$$

since the number of times $e_{-\mu}$ appears on the right-hand side is the number of sets $\left(n_{\alpha, i}\right)$ of non-negative integers such that

$$
\sum_{\alpha \in \Phi^{+}}\left(\sum_{i=1}^{m_{\alpha}} n_{\alpha, i}\right) \alpha=\mu
$$

Hence the character of $M(\lambda)$ is

$$
\operatorname{ch} M(\lambda)=e_{\lambda} \prod_{\alpha \in \Phi^{+}}\left(1+e_{-\alpha}+e_{-\alpha}^{2}+\cdots\right)^{m_{\alpha}}
$$

Now the element $1+e_{-\alpha}+e_{-\alpha}^{2}+\cdots \in \mathfrak{R}$ has inverse $1-e_{-\alpha} \in \mathfrak{R}$. Thus we have

$$
\operatorname{ch} M(\lambda)=\frac{e_{\lambda}}{\prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)^{m_{\alpha}}}
$$

Of course in the special case when $L(A)$ is finite dimensional this formula reduces to that obtained in Proposition 12.4. In the general case there are two differences - the roots need not have multiplicity 1 and the product over the positive roots can be an infinite product.

Now the Verma module $M(\lambda)$ for $L(A)$ has a unique maximal submodule $J(\lambda)$, just as in the proof of Theorem 10.9. We define

$$
L(\lambda)=M(\lambda) / J(\lambda)
$$

Then $L(\lambda)$ is an irreducible $L(A)$-module in the category $\mathcal{O}$.

Proposition 19.3 The modules $L(\lambda)$ for $\lambda \in H^{*}$ are the only irreducible modules in category $\mathcal{O}$.

Proof. Let $V$ be an irreducible $L(A)$-module with $V \in \mathcal{O}$. The definition of $\mathcal{O}$ shows that $V$ has a maximal weight $\lambda$ under the partial ordering $\prec$. Let $v_{\lambda} \in V$ be a weight vector with weight $\lambda$. Then $x v_{\lambda}=0$ for $x \in N$ and $x v_{\lambda}=\lambda(x) v_{\lambda}$ for $x \in H$.

We may define a map from the Verma module $M(\lambda)$ into $V$ as follows. Each element of $M(\lambda)$ has a unique expression of form $u m_{\lambda}$ for $u \in \mathfrak{U}\left(N^{-}\right)$. Let $\theta: M(\lambda) \rightarrow V$ be defined by $\theta\left(u m_{\lambda}\right)=u v_{\lambda}$ for $u \in \mathfrak{U}\left(N^{-}\right)$. It may then be shown, just as in the proof of Proposition 10.13, that $\theta$ is a homomorphism of $\mathfrak{H}(L(A))$-modules. The image of $\theta$ is a submodule of $V$ containing $v_{\lambda}$, so is the whole of $V$ since $V$ is irreducible. Thus the kernel of $\theta$ is a maximal submodule of $M(\lambda)$, so must be $J(\lambda)$. Thus $V$ is isomorphic to $M(\lambda) / J(\lambda)=L(\lambda)$.

Now in Theorem 12.16 we showed that each Verma module $M(\lambda)$ for a finite dimensional semisimple Lie algebra $L$ has a finite composition series. The proof of this result made extensive use of the fact that $L$ is finite dimensional, and the result does not carry over to Verma modules for Kac-Moody algebras $L(A)$. For example it can be shown that the Verma module $M(0)$ has no irreducible submodule when $L(A)$ is infinite dimensional.

We would nevertheless like to define the multiplicity $[V: L(\lambda)]$ of the irreducible module $L(\lambda)$ in the module $V \in \mathcal{O}$. If $V$ had a finite composition series $[V: L(\lambda)$ ] would be the number of composition factors isomorphic to $L(\lambda)$ in a given composition series, and this would be independent of the choice of composition series by the Jordan-Hölder theorem. However, $V$ does not in general have a finite composition series. Even so, Kac found a way of defining the multiplicity $[V: L(\lambda)]$. This makes use of the following lemma.

Lemma 19.4 Let $V \in \mathcal{O}$ and $\lambda \in H^{*}$. Then $V$ has a filtration

$$
V=V_{0} \supset V_{1} \supset \cdots \supset V_{t}=O
$$

of finite length by means of a sequence of submodules such that each factor $V_{i-1} / V_{i}$ either is isomorphic to $L(\mu)$ for some $\mu \succ \lambda$ or has the property that $\left(V_{i-1} / V_{i}\right)_{\mu}=O$ for all $\mu \succ \lambda$.

Proof. The definition of $\mathcal{O}$ shows that $V$ has only finitely many weights $\mu$ with $\mu \succ \lambda$. Thus

$$
a(V, \lambda)=\sum_{\mu \succ \lambda} \operatorname{dim} V_{\mu}
$$

is finite. We shall prove the lemma by induction on $a(V, \lambda)$. If $a(V, \lambda)=0$ then $V=V_{0} \supset V_{1}=O$ is the required filtration. So suppose $a(V, \lambda)>0$. Then $V$ has a weight $\mu$ with $\mu \succ \lambda$. We may choose a maximal weight $\mu$ with $\mu \succ \lambda$. Let $v_{\mu} \in V$ be a weight vector with weight $\mu$. Then $x v_{\mu}=0$ for $x \in N$ and $x v_{\mu}=\mu(x) \nu_{\mu}$ for $x \in H$. Let $U=\mathfrak{l l}(L(A)) v_{\mu}$ be the submodule of $V$
generated by $v_{\mu}$. We then have a map $\theta: M(\mu) \rightarrow U$ defined by $\theta\left(u m_{\mu}\right)=$ $u v_{\mu}$ for $u \in \mathfrak{H}\left(N^{-}\right)$, and $\theta$ is a homomorphism of $\mathfrak{H}(L(A))$-modules as before, as shown in the proof of Proposition 10.13. Moreover $\theta$ is surjective. Thus $U$ is isomorphic to a factor module of the Verma module $M(\mu)$, and so has a unique maximal submodule $\bar{U}$. We also have

$$
U / \bar{U} \cong M(\mu) / J(\mu) \cong L(\mu)
$$

Now consider the filtration

$$
V \supset U \supset \bar{U} \supset O
$$

We have $a(\bar{U}, \lambda)<a(V, \lambda)$ and $a(V / U, \lambda)<a(V, \lambda)$, since the weight $\mu \succ \lambda$ appears in $U / \bar{U}$. Thus by induction we obtain filtrations for the modules $\bar{U} \in \mathcal{O}$ and $V / U \in \mathcal{O}$ of the required kind, and these may be combined to give the required filtration of $V$.

Lemma 19.5 Let $V \in \mathcal{O}$ and $\lambda \in H^{*}$. Consider filtrations of the type given in Lemma 19.4 with respect to $\lambda$. Let $\mu \in H^{*}$ satisfy $\mu \succ \lambda$. Then the number of factors $L(\mu)$ in such a filtration is independent of the choice of filtration and also of the choice of $\lambda$.

Proof. We first observe that a filtration with respect to $\lambda$ is also a filtration with respect to $\mu$ when $\mu \succ \lambda$. Also the multiplicity of $L(\mu)$ in such a filtration is the same whether it is regarded as a filtration with respect to $\lambda$ or $\mu$. Thus to prove the lemma it will be sufficient to take two filtrations with respect to $\mu$ and show that $L(\mu)$ has the same multiplicity in each.

The following variant of the proof of the Jordan-Hölder theorem achieves this. Let

$$
\begin{align*}
& V=V_{0} \supset V_{1} \supset \cdots \supset V_{l_{1}}=O  \tag{19.1}\\
& V=V_{0}^{\prime} \supset V_{1}^{\prime} \supset \cdots \supset V_{l_{2}}^{\prime}=O \tag{19.2}
\end{align*}
$$

be two such filtrations of lengths $l_{1}, l_{2}$. We shall use induction on $\min \left(l_{1}, l_{2}\right)$.
Suppose first that $\min \left(l_{1}, l_{2}\right)=1$. Then either $V$ is irreducible and the two filtrations are identical, or $\mu$ is not a weight of $V$ and $L(\mu)$ does not appear in either filtration.

Thus suppose $\min \left(l_{1}, l_{2}\right)>1$. We suppose first that $V_{1}=V_{1}^{\prime}$. We then consider the two filtrations

$$
\begin{aligned}
& V_{1} \supset \cdots \supset V_{l_{1}}=O \\
& V_{1}^{\prime} \supset \cdots \supset V_{l_{2}}^{\prime}=O
\end{aligned}
$$

of $V_{1}$. By induction they give the same multiplicity for $L(\mu)$, and the filtrations for $V$ are obtained by adding the additional factor $V / V_{1}$ which is the same for both.

We may therefore suppose that $V_{1} \neq V_{1}^{\prime}$. Suppose first that one contains the other, say $V_{1} \subset V_{1}^{\prime}$. Then $V / V_{1}$ is not irreducible and so $\mu$ is not a weight of $V / V_{1}$. Thus neither $V / V_{1}$ nor $V / V_{1}^{\prime}$ is isomorphic to $L(\mu)$. Let

$$
V_{1} \supset U_{1} \supset \cdots \supset U_{m}=O
$$

be a filtration of $V_{1}$ of the required type with respect to $\mu$. We then consider the filtrations

$$
\begin{array}{r}
V \supset V_{1} \supset U_{1} \supset \cdots \supset U_{m}=O \\
V \supset V_{1}^{\prime} \supset V_{1} \supset U_{1} \supset \cdots \supset U_{m}=O . \tag{19.4}
\end{array}
$$

These are filtrations of $V$ of the required type with respect to $\mu . L(\mu)$ has the same multiplicity in filtrations (19.1), (19.3) since they have the same leading term $V_{1}$. Similarly $L(\mu)$ has the same multiplicity in filtrations (19.2), (19.4). So $L(\mu)$ has the same multiplicity in filtrations (19.3), (19.4) since none of $V / V_{1}, V / V_{1}^{\prime}, V_{1}^{\prime} / V_{1}$ is isomorphic to $L(\mu)$. Thus $L(\mu)$ has the same multiplicity in filtrations (19.1), (19.2) as required.

We may therefore assume that neither of $V_{1}, V_{1}^{\prime}$ is contained in the other. Let $U=V_{1} \cap V_{1}^{\prime}$ and choose a filtration of $U$ of the required kind with respect to $\mu$. This has form

$$
U \supset U_{1} \supset \cdots \supset U_{m}=O
$$

We then consider the filtrations

$$
\begin{gather*}
V \supset V_{1} \supset U \supset U_{1} \supset \cdots \supset U_{m}=O  \tag{19.5}\\
V \supset V_{1}^{\prime} \supset U \supset U_{1} \supset \cdots \supset U_{m}=O \tag{19.6}
\end{gather*}
$$

These are filtrations of $V$ of the required type with respect to $\mu$. This is clear since

$$
V_{1} / U \cong\left(V_{1}+V_{1}^{\prime}\right) / V_{1}^{\prime}, \quad V_{1}^{\prime} / U \cong\left(V_{1}+V_{1}^{\prime}\right) / V_{1}
$$

Now $L(\mu)$ has the same multiplicity in filtrations (19.1), (19.5) and the same multiplicity in filtrations (19.2), (19.6) since the leading terms are the same. It is therefore sufficient to show that $L(\mu)$ has the same multiplicity in filtrations (19.5), (19.6). These filtrations differ only in the first two factors. If $V_{1}+V_{1}^{\prime}=V$ then we have

$$
V / V_{1} \cong V_{1}^{\prime} / U, \quad V / V_{1}^{\prime} \cong V_{1} / U
$$

as required. If $V_{1}+V_{1}^{\prime} \neq V$ then $V / V_{1}$ and $V / V_{1}^{\prime}$ are not irreducible. In this case $\mu$ is not a weight of $V / V_{1}$ or $V / V_{1}^{\prime}$, so is not a weight of $V_{1} / U$. Thus none of $V / V_{1}, V_{1} / U, V / V_{1}^{\prime}, V_{1}^{\prime} / U$ is isomorphic to $L(\mu)$. This completes the proof.

Definition The multiplicity of $L(\mu)$ in a filtration of $V \in \mathcal{O}$ of the type considered in Lemmas 19.4 and 19.5 will be denoted by $[V: L(\mu)]$.

Of course this agrees with the previous definition of $[V: L(\mu)]$ in the case when $V$ has a composition series of finite length.

Proposition 19.6 Let $V \in \mathcal{O}$. Then

$$
\operatorname{ch} V=\sum_{\lambda \in H^{*}}[V: L(\lambda)] \operatorname{ch} L(\lambda)
$$

Proof. Both sides are functions $H^{*} \rightarrow \mathbb{Z}$. We have $(\operatorname{ch} V)(\mu)=\operatorname{dim} V_{\mu}$ and the right-hand side evaluated at $\mu$ is

$$
\sum_{\lambda \in H^{*}}[V: L(\lambda)] \operatorname{dim} L(\lambda)_{\mu} .
$$

We choose a filtration of $V$ with respect to $\mu$ of the type given in Lemma 19.4. Each factor either is isomorphic to $L(\lambda)$ for some $\lambda \succ \mu$ or does not contain $\mu$ as a weight. The multiplicity of $L(\lambda)$ as a factor is $[V: L(\lambda)]$. Hence we have

$$
\operatorname{dim} V_{\mu}=\sum_{\lambda}[V: L(\lambda)] \operatorname{dim} L(\lambda)_{\mu}
$$

summed over all $\lambda \succ \mu$. We may in fact take the sum over all $\lambda \in H^{*}$ since $\operatorname{dim} L(\lambda)_{\mu}=0$ unless $\lambda \succ \mu$.

### 19.2 The generalised Casimir operator

We recall from Section 11.6 that, if $L$ is a finite dimensional semisimple Lie algebra, the Casimir element of the centre of the enveloping algebra $\mathfrak{H}(L)$ plays an important role in the representation theory of $L$. If $x_{1}, \ldots, x_{m}$ are any basis of $L$ and $y_{1}, \ldots, y_{m}$ are the dual basis with respect to the Killing form the Casimir element is given by

$$
\sum x_{i} y_{i} \in \mathfrak{H}(L)
$$

We showed in Proposition 11.36 that the Casimir element acts on a Verma module $M(\lambda)$ for $L$ as scalar multiplication by $\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle$ where $\langle$,$\rangle is the Killing form and \rho$ is, as usual, the element of $H^{*}$ given by $\rho\left(h_{i}\right)=1$ for $i=1, \ldots, l$.

Now let $L(A)$ be a Kac-Moody algebra where $A$ is symmetrisable. We cannot define an analogous Casimir element $\sum x_{i} y_{i}$ in $\mathfrak{U}(L(A))$ since the sum will in general be infinite and make no sense. It was shown by Kac, however, that it is possible to define an operator $c: V \rightarrow V$ on any $L(A)$-module $V$ in category $\mathcal{O}$ which has properties analogous to the action of the Casimir element for finite dimensional algebras. In order to define Kac' operator on $V$ we recall the formula for the Casimir element of a finite dimensional algebra given in Proposition 11.35. Let $h_{1}^{\prime}, \ldots, h_{l}^{\prime}$ be a basis of $H$ and $h_{1}^{\prime \prime}, \ldots, h_{l}^{\prime \prime}$ be the dual basis of $H$ with respect to the Killing form of $L$. Choose elements $e_{\alpha} \in L_{\alpha}, f_{\alpha} \in L_{-\alpha}$ such that $\left[e_{\alpha} f_{\alpha}\right]=h_{\alpha}^{\prime}$ for each $\alpha \in \Phi^{+}$. Then the Casimir element of $\mathfrak{H}(L)$ is given by

$$
\sum_{i=1}^{l} h_{i}^{\prime} h_{i}^{\prime \prime}+\sum_{\alpha \in \Phi^{+}} h_{\alpha}^{\prime}+2 \sum_{\alpha \in \Phi^{+}} f_{\alpha} e_{\alpha} .
$$

Since $\sum_{\alpha \in \Phi^{+}} \alpha=2 \rho$ this element can also be written

$$
\sum_{i=1}^{l} h_{i}^{\prime} h_{i}^{\prime \prime}+2 h_{\rho}^{\prime}+2 \sum_{\alpha \in \Phi^{+}} f_{\alpha} e_{\alpha}
$$

where $h_{\rho}^{\prime} \in H$ satisfies $\rho(x)=\left\langle h_{\rho}^{\prime}, x\right\rangle$ for all $x \in H$.
We wish to define an analogous element for the symmetrisable Kac-Moody algebra $L(A)$. The root space $L_{\alpha}$ of $L(A)$ need not be 1-dimensional, so we choose a basis $e_{\alpha}^{(1)}, e_{\alpha}^{(2)}, \ldots$ for $L_{\alpha}$. Instead of using the Killing form we use the standard invariant bilinear form on $L(A)$. (In the case when $L(A)$ is finite dimensional this is a scalar multiple of the Killing form.) We recall from Corollary 16.5 that the pairing $L_{\alpha} \times L_{-\alpha} \rightarrow \mathbb{C}$ given by $x, y \rightarrow\langle x, y\rangle$ is non-degenerate. Thus we may choose a corresponding dual basis $f_{\alpha}^{(1)}, f_{\alpha}^{(2)}, \ldots$ for $L_{-\alpha}$ such that

$$
\left\langle e_{\alpha}^{(i)}, f_{\alpha}^{(j)}\right\rangle=\delta_{i j} .
$$

We choose a basis $h_{1}^{\prime}, h_{2}^{\prime}, \ldots$ of $H$ and let $h_{1}^{\prime \prime}, h_{2}^{\prime \prime}, \ldots$ be the dual basis of $H$ satisfying $\left\langle h_{i}^{\prime}, h_{j}^{\prime \prime}\right\rangle=\delta_{i j}$. Since the fundamental coroots $h_{1}, \ldots, h_{n} \in H$ are linearly independent there exists $\rho \in H^{*}$ such that $\rho\left(h_{i}\right)=1$ for $i=1, \ldots, n$. However, $\rho$ is not in general uniquely determined by this condition. So we choose any element $\rho \in H^{*}$ satisfying $\rho\left(h_{i}\right)=1$ for $i=1, \ldots, n$. We then have a corresponding element $h_{\rho}^{\prime} \in H$ such that $\rho(x)=\left\langle h_{\rho}^{\prime}, x\right\rangle$ for all $x \in H$. We then consider the expression

$$
\sum_{i} h_{i}^{\prime} h_{i}^{\prime \prime}+2 h_{\rho}^{\prime}+2 \sum_{\alpha \in \Phi^{+}} \sum_{i} f_{\alpha}^{(i)} e_{\alpha}^{(i)} .
$$

This element does not make sense as an element of $\mathfrak{l}(L(A))$ in general since the sum over $\alpha \in \Phi^{+}$may be infinite. However, if $V$ is an $L(A)$-module in
$\mathcal{O}$ we know that ch $V \in \mathfrak{R}$ and so there exist only finitely many $\alpha \in \Phi^{+}$such that $L_{\alpha} V \neq O$. Thus the operator $\Omega: V \rightarrow V$ given by

$$
\Omega=\sum_{i} h_{i}^{\prime} h_{i}^{\prime \prime}+2 h_{\rho}^{\prime}+2 \sum_{\alpha \in \Phi^{+}} \sum_{i} f_{\alpha}^{(i)} e_{\alpha}^{(i)}
$$

is well defined. It is straightforward to check that this operator $\Omega: V \rightarrow V$ does not depend on the choice of dual bases $h_{1}^{\prime}, h_{2}^{\prime}, \ldots \quad h_{1}^{\prime \prime}, h_{2}^{\prime \prime}, \ldots$ of $H$ or on the choice of dual bases $e_{\alpha}^{(i)}, f_{\alpha}^{(i)}$ for $L_{\alpha}$ and $L_{-\alpha}$. It may, however, depend upon the choice of $\rho$.

Definition The operator $\Omega: V \rightarrow V$ for $V \in \mathcal{O}$ is called the generalised Casimir operator on $V$ with respect to $\rho$.

In the case of a finite dimensional semisimple Lie algebra the Casimir element lies in the centre of the universal enveloping algebra. We shall prove an analogous result in the present situation, i.e. that the generalised Casimir operator commutes with the action on $V \in \mathcal{O}$ of any element of $\mathfrak{H}(L(A))$. We first need some preliminary results.

Lemma 19.7 Let $\alpha, \beta, \beta-\alpha \in \Phi$. Suppose $e_{\alpha}^{(i)}, f_{\alpha}^{(i)}$ are dual bases of $L_{\alpha}, L_{-\alpha}$ and $e_{\beta}^{(i)}, f_{\beta}^{(i)}$ are dual bases of $L_{\beta}, L_{-\beta}$. Let $x \in L_{\beta-\alpha}$. Then in the vector space $L(A) \otimes L(A)$ we have

$$
\sum_{i} f_{\alpha}^{(i)} \otimes\left[x, e_{\alpha}^{(i)}\right]=\sum_{i}\left[f_{\beta}^{(i)}, x\right] \otimes e_{\beta}^{(i)}
$$

Proof. We note that both sides lie in the subspace $L_{-\alpha} \otimes L_{\beta}$. We define a bilinear form on $L(A) \otimes L(A)$, uniquely determined by

$$
\left\langle x_{1} \otimes y_{1}, x_{2} \otimes y_{2}\right\rangle=\left\langle x_{1}, x_{2}\right\rangle\left\langle y_{1}, y_{2}\right\rangle .
$$

Since the standard invariant form is non-degenerate on $L(A)$ this bilinear form will be non-degenerate on $L(A) \otimes L(A)$.

Let $a \otimes b \in L_{\gamma} \otimes L_{\delta}$. Then the scalar products of both sides of the required equation with $a \otimes b$ are zero unless $\gamma=\alpha$ and $\delta=-\beta$. We therefore suppose $\gamma=\alpha$ and $\delta=-\beta$ and consider the scalar products

$$
\begin{aligned}
& \left\langle\sum_{i} f_{\alpha}^{(i)} \otimes\left[x, e_{\alpha}^{(i)}\right], a \otimes b\right\rangle \\
& \left\langle\sum_{i}\left[f_{\beta}^{(i)}, x\right] \otimes e_{\beta}^{(i)}, a \otimes b\right\rangle
\end{aligned}
$$

We have

$$
\begin{aligned}
\left\langle\sum_{i} f_{\alpha}^{(i)} \otimes\left[x, e_{\alpha}^{(i)}\right], a \otimes b\right\rangle & =\sum_{i}\left\langle f_{\alpha}^{(i)}, a\right\rangle\left\langle\left[x, e_{\alpha}^{(i)}\right], b\right\rangle \\
& =-\sum_{i}\left\langle f_{\alpha}^{(i)}, a\right\rangle\left\langle e_{\alpha}^{(i)},[x b]\right\rangle \\
& =-\langle a,[x b]\rangle
\end{aligned}
$$

since $e_{\alpha}^{(i)}, f_{\alpha}^{(i)}$ are dual bases of $L_{\alpha}, L_{-\alpha}$. Similarly

$$
\begin{aligned}
\left\langle\sum_{i}\left[f_{\beta}^{(i)}, x\right] \otimes e_{\beta}^{(i)}, a \otimes b\right\rangle & =\sum_{i}\left\langle\left[f_{\beta}^{(i)}, x\right], a\right\rangle\left\langle e_{\beta}^{(i)}, b\right\rangle \\
& =\sum_{i}\left\langle f_{\beta}^{(i)},[x, a]\right\rangle\left\langle e_{\beta}^{(i)}, b\right\rangle \\
& =\langle[x a], b\rangle \\
& =-\langle a,[x b]\rangle .
\end{aligned}
$$

Thus the two sides of our equation have the same scalar product with each $a \otimes$ $b \in L_{\alpha} \otimes L_{-\beta}$. Since the form is non-degenerate on $L(A) \otimes L(A)$ this shows the two sides are equal.

Corollary 19.8 In the enveloping algebra $\mathfrak{H}(L(A))$ we have

$$
\sum_{i} f_{\alpha}^{(i)}\left[x, e_{\alpha}^{(i)}\right]=\sum_{i}\left[f_{\beta}^{(i)}, x\right] e_{\beta}^{(i)} .
$$

Proof. We apply the natural homomorphism from the tensor algebra $T(L(A))$ to $\mathfrak{H}(L(A))$. The result then follows from Lemma 19.7.

Theorem 19.9 Let $u \in \mathfrak{H}(L(A))$ and $V \in \mathcal{O}$. Then the maps $\Omega: V \rightarrow V$ and $u: V \rightarrow V$ commute.

Proof. The algebra $\mathfrak{U}(L(A))$ is generated by $e_{i}, f_{i}$ for $i=1, \ldots, n$ and the elements of $H$. If $x \in H$ then $x$ commutes with each term $f_{\alpha}^{(i)} e_{\alpha}^{(i)}$ in $\mathfrak{H}(L(A))$ since this term has weight 0 . Thus $x: V \rightarrow V$ commutes with $f_{\alpha}^{(i)} e_{\alpha}^{(i)}: V \rightarrow V$ and hence with $\Omega: V \rightarrow V$. It is therefore sufficient to show that $\Omega: V \rightarrow V$ commutes with $e_{i}: V \rightarrow V$ and $f_{i}: V \rightarrow V$.

We consider the element $\left[\sum_{j} f_{\alpha}^{(j)} e_{\alpha}^{(j)}, e_{i}\right]$ of $\mathfrak{l}(L(A))$. We have

$$
\begin{aligned}
{\left[\sum_{j} f_{\alpha}^{(j)} e_{\alpha}^{(j)}, e_{i}\right] } & =\sum_{j}\left[f_{\alpha}^{(j)}, e_{i}\right] e_{\alpha}^{(j)}+\sum_{j} f_{\alpha}^{(j)}\left[e_{\alpha}^{(j)}, e_{i}\right] \\
& =\sum_{j}\left[f_{\alpha}^{(j)}, e_{i}\right] e_{\alpha}^{(j)}-\sum_{j} f_{\alpha}^{(j)}\left[e_{i}, e_{\alpha}^{(j)}\right] \\
& =\sum_{j}\left[f_{\alpha}^{(j)}, e_{i}\right] e_{\alpha}^{(j)}-\sum_{j}\left[f_{\alpha+\alpha_{i}}^{(j)}, e_{i}\right] e_{\alpha+\alpha_{i}}^{(j)}
\end{aligned}
$$

by Corollary 19.8. If $\alpha+\alpha_{i} \notin \Phi$ the second term is interpreted as 0 .
We show that

$$
\sum_{\substack{\alpha \in \Phi^{+} \\ \alpha \neq \alpha_{i}}}\left(\sum_{j} f_{\alpha}^{(j)} e_{\alpha}^{(j)}\right): V \rightarrow V
$$

commutes with $e_{i}: V \rightarrow V$. We have

$$
\left[\sum_{\substack{\alpha \in \Phi^{+} \\ \alpha \neq \alpha_{i}}}\left(\sum_{j} f_{\alpha}^{(j)} e_{\alpha}^{(j)}\right), e_{i}\right]=\sum_{\substack{\alpha \in \Phi^{+} \\ \alpha \neq \alpha_{i}}}\left[\sum_{j} f_{\alpha}^{(j)}, e_{i}\right] e_{\alpha}^{(j)}-\sum_{\substack{\alpha \in \Phi^{+} \\ \alpha \neq \alpha_{i}}} \sum_{j}\left[f_{\alpha+\alpha_{i}}^{(j)}, e_{i}\right] e_{\alpha+\alpha_{i}}^{(j)}
$$

on $V$. If $\alpha-\alpha_{i} \notin \Phi$ then $\left[\sum_{j} f_{\alpha}^{(j)}, e_{i}\right]=0$. Thus we may assume $\alpha=\beta+\alpha_{i}$ in the first term with $\beta \in \Phi^{+}$and get

$$
\sum_{\substack{\beta \in \Phi^{+} \\ \beta \neq \alpha_{i}}}\left[\sum_{j} f_{\beta+\alpha_{i}}^{(j)}, e_{i}\right] e_{\beta+\alpha_{i}}^{(j)}-\sum_{\substack{\alpha \in \Phi^{+} \\ \alpha \neq \alpha_{i}}} \sum_{j}\left[f_{\alpha+\alpha_{i}}^{(j)}, e_{i}\right] e_{\alpha+\alpha_{i}}^{(j)}=0 .
$$

Since $\Omega=\sum_{j} h_{j}^{\prime} h_{j}^{\prime \prime}+2 h_{\rho}^{\prime}+2 \sum_{\alpha \in \Phi^{+}} \sum_{j} f_{\alpha}^{(j)} e_{\alpha}^{(j)}$ on $V$ it is now sufficient to show that $\sum_{j} h_{j}^{\prime} h_{j}^{\prime \prime}+2 h_{\rho}^{\prime}+2 f_{i} e_{i}$ commutes with $e_{i}$ on $V$. In fact these elements commute in $\mathfrak{U}(L(A))$. For we have

$$
\begin{aligned}
{\left[\sum_{j} h_{j}^{\prime} h_{j}^{\prime \prime}, e_{i}\right] } & =\sum_{j}\left[h_{j}^{\prime}, e_{i}\right] h_{j}^{\prime \prime}+\sum_{j} h_{j}^{\prime}\left[h_{j}^{\prime \prime}, e_{i}\right] \\
& =\sum_{j} \alpha_{i}\left(h_{j}^{\prime}\right) e_{i} h_{j}^{\prime \prime}+\sum_{j} \alpha_{i}\left(h_{j}^{\prime \prime}\right) h_{j}^{\prime} e_{i} \\
& =e_{i}\left(\sum_{j} \alpha_{i}\left(h_{j}^{\prime}\right) h_{j}^{\prime \prime}+\sum_{j} \alpha_{i}\left(h_{j}^{\prime \prime}\right) h_{j}^{\prime}\right)+\left(\sum_{j} \alpha_{i}\left(h_{j}^{\prime \prime}\right) \alpha_{i}\left(h_{j}^{\prime}\right)\right) e_{i} \\
& =e_{i}\left(\sum_{j}\left\langle h_{\alpha_{i}}^{\prime}, h_{j}^{\prime}\right\rangle h_{j}^{\prime \prime}+\sum_{j}\left\langle h_{\alpha_{i}}^{\prime}, h_{j}^{\prime \prime}\right\rangle h_{j}^{\prime}\right)+\left(\sum_{j}\left\langle h_{\alpha_{i}}^{\prime}, h_{j}^{\prime \prime}\right\rangle\left\langle h_{\alpha_{i}}^{\prime}, h_{j}^{\prime}\right\rangle\right) e_{i}
\end{aligned}
$$

Since $h_{1}^{\prime}, h_{2}^{\prime}, \ldots$ and $h_{1}^{\prime \prime}, h_{2}^{\prime \prime}, \ldots$ are dual bases of $H$ we have

$$
\sum_{j}\left\langle h_{\alpha_{i}}^{\prime}, h_{j}^{\prime}\right\rangle h_{j}^{\prime \prime}=\sum_{j}\left\langle h_{\alpha_{i}}^{\prime}, h_{j}^{\prime \prime}\right\rangle h_{j}^{\prime}=h_{\alpha_{i}}^{\prime}
$$

and

$$
\sum_{j}\left\langle h_{\alpha_{i}}^{\prime}, h_{j}^{\prime \prime}\right\rangle\left\langle h_{\alpha_{i}}^{\prime}, h_{j}^{\prime}\right\rangle=\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle=\left\langle\alpha_{i}, \alpha_{i}\right\rangle .
$$

Hence

$$
\left[\sum_{j} h_{j}^{\prime} h_{j}^{\prime \prime}, e_{i}\right]=2 e_{i} h_{\alpha_{i}}^{\prime}+\left\langle\alpha_{i}, \alpha_{i}\right\rangle e_{i} .
$$

Secondly we have

$$
\left[2 h_{\rho}^{\prime}, e_{i}\right]=2 \alpha_{i}\left(h_{\rho}^{\prime}\right) e_{i}=2\left\langle h_{\alpha_{i}}^{\prime}, h_{\rho}^{\prime}\right\rangle e_{i}=2 \rho\left(h_{\alpha_{i}}^{\prime}\right) e_{i}=\left\langle\alpha_{i}, \alpha_{i}\right\rangle e_{i}
$$

since $h_{i}=\frac{2 h_{\alpha_{i}}^{\prime}}{\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{i}}^{\prime}\right|}$ and $\rho\left(h_{i}\right)=1$, hence

$$
\rho\left(h_{\alpha_{i}}^{\prime}\right)=\frac{\left\langle h_{\alpha_{i}}^{\prime}, h_{\alpha_{i}}^{\prime}\right\rangle}{2}=\frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{2} .
$$

Thirdly we have

$$
\left[2 f_{i} e_{i}, e_{i}\right]=2\left[f_{i}, e_{i}\right] e_{i}=-2\left[e_{i}, f_{i}\right] e_{i}
$$

We recall that $e_{i}, f_{i}$ were chosen so that $\left\langle e_{i}, f_{i}\right\rangle=1$. By Corollary 16.5 this implies $\left[e_{i} f_{i}\right]=h_{\alpha_{i}}^{\prime}$. Hence

$$
\begin{aligned}
{\left[2 f_{i} e_{i}, e_{i}\right] } & =-2 h_{\alpha_{i}}^{\prime} e_{i}=-2 e_{i} h_{\alpha_{i}}^{\prime}-2 \alpha_{i}\left(h_{\alpha_{i}}^{\prime}\right) e_{i} \\
& =-2 e_{i} h_{\alpha_{i}}^{\prime}-2\left\langle\alpha_{i}, \alpha_{i}\right\rangle e_{i} .
\end{aligned}
$$

Thus we have shown:

$$
\begin{aligned}
{\left[\sum_{j} h_{j}^{\prime} h_{j}^{\prime \prime}, e_{i}\right] } & =2 e_{i} h_{\alpha_{i}}^{\prime}+\left\langle\alpha_{i}, \alpha_{i}\right\rangle e_{i} \\
{\left[2 h_{\rho}^{\prime}, e_{i}\right] } & =\left\langle\alpha_{i}, \alpha_{i}\right\rangle e_{i} \\
{\left[2 f_{i} e_{i}, e_{i}\right] } & =-2 e_{i} h_{\alpha_{i}}^{\prime}-2\left\langle\alpha_{i}, \alpha_{i}\right\rangle e_{i} .
\end{aligned}
$$

Hence

$$
\left[\sum_{j} h_{j}^{\prime} h_{j}^{\prime \prime}+2 h_{\rho}^{\prime}+2 f_{i} e_{i}, e_{i}\right]=0 .
$$

Thus we have shown that $\Omega: V \rightarrow V$ commutes with $e_{i}: V \rightarrow V$. The proof that $\Omega: V \rightarrow V$ commutes with $f_{i}: V \rightarrow V$ is similar. Using the fact that

$$
\begin{aligned}
{\left[\sum_{j} f_{\alpha}^{(j)} e_{\alpha}^{(j)}, f_{i}\right] } & =\sum_{j}\left[f_{\alpha}^{(j)}, f_{i}\right] e_{\alpha}^{(j)}+\sum_{j} f_{\alpha}^{(j)}\left[e_{\alpha}^{(j)}, f_{i}\right] \\
& =\sum_{j} f_{\alpha+\alpha_{i}}^{(j)}\left[f_{i}, e_{\alpha+\alpha_{i}}^{(j)}\right]-\sum_{j} f_{\alpha}^{(j)}\left[f_{i}, e_{\alpha}^{(j)}\right]
\end{aligned}
$$

we deduce as before that

$$
\sum_{\substack{\alpha \in \Phi^{+} \\ \alpha \neq \alpha_{i}}}\left(\sum_{j} f_{\alpha}^{(j)} e_{\alpha}^{(j)}\right): V \rightarrow V
$$

commutes with $f_{i}: V \rightarrow V$. We also obtain

$$
\begin{aligned}
{\left[\sum_{j} h_{j}^{\prime} h_{j}^{\prime \prime}, f_{i}\right] } & =-2 f_{i} h_{\alpha_{i}}^{\prime}+\left\langle\alpha_{i}, \alpha_{i}\right\rangle f_{i} \\
{\left[2 h_{\rho}^{\prime}, f_{i}\right] } & =-\left\langle\alpha_{i}, \alpha_{i}\right\rangle f_{i} \\
{\left[2 f_{i} e_{i}, f_{i}\right] } & =2 f_{i} h_{\alpha_{i}}^{\prime}
\end{aligned}
$$

Hence

$$
\left[\sum_{j} h_{j}^{\prime} h_{j}^{\prime \prime}+2 h_{\rho}^{\prime}+2 f_{i} e_{i}, f_{i}\right]=0
$$

Thus $\Omega: V \rightarrow V$ commutes with $f_{i}: V \rightarrow V$ and the proof is complete.
We next describe the action of the generalised Casimir operator $\Omega$ on a Verma module.

Proposition 19.10 $\Omega$ acts on the Verma module $M(\lambda)$ as scalar multiplication by $\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle$.

Proof. Let $m_{\lambda}$ be a highest weight vector of $M(\lambda)$. Then

$$
\begin{aligned}
\Omega m_{\lambda} & =\left(\sum_{j} h_{j}^{\prime} h_{j}^{\prime \prime}+2 h_{\rho}^{\prime}+2 \sum_{\alpha \in \Phi^{+}} \sum_{j} f_{\alpha}^{(j)} e_{\alpha}^{(j)}\right) m_{\lambda} \\
& =\left(\sum_{j} \lambda\left(h_{j}^{\prime}\right) \lambda\left(h_{j}^{\prime \prime}\right)+2 \lambda\left(h_{\rho}^{\prime}\right)\right) m_{\lambda}
\end{aligned}
$$

Now $\sum_{j} \lambda\left(h_{j}^{\prime}\right) \lambda\left(h_{j}^{\prime \prime}\right)=\langle\lambda, \lambda\rangle$ and $\lambda\left(h_{\rho}^{\prime}\right)=\langle\lambda, \rho\rangle$. Hence

$$
\Omega m_{\lambda}=(\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle) m_{\lambda} .
$$

Now each element of $M(\lambda)$ has form $u m_{\lambda}$ for some $u \in \mathfrak{H}\left(N^{-}\right)$. Thus

$$
\Omega\left(u m_{\lambda}\right)=u\left(\Omega m_{\lambda}\right)=(\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle) u m_{\lambda},
$$

by Theorem 19.9. Hence $\Omega$ acts on $M(\lambda)$ as scalar multiplication by

$$
\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle .
$$

Corollary $19.11 \Omega$ acts on the irreducible $L(A)$-module $L(\lambda)$ as scalar multiplication by $\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle$

Note Proposition 19.10 is the analogue of Proposition 11.36 for finite dimensional semisimple Lie algebras. In Proposition 11.36 the invariant form which appeared was the Killing form whereas in Proposition 19.10 and Corollary 19.11 it is the standard invariant form. The difference is explained by the fact that the Casimir element in the enveloping algebra of a finite dimensional semisimple Lie algebra was defined in terms of the Killing form, whereas the generalised Casimir operator was defined in terms of the standard invariant form.

### 19.3 Kac' character formula

Let $X$ be the set of integral weights $\lambda \in H^{*}$, that is the set of all $\lambda$ such that $\lambda\left(h_{i}\right) \in \mathbb{Z}$ for $i=1, \ldots, n$. Let $X^{+}$be the subset of dominant integral weights, that is the set of weights $\lambda \in X$ such that $\lambda\left(h_{i}\right) \geq 0$ for all $i$. In this section we shall prove a formula due to Kac for the character of the irreducible $L(A)$-module $L(\lambda)$ when $\lambda \in X^{+}$. The reason for the restriction to weights in $X^{+}$lies in the fact that the modules $L(\lambda)$ for $\lambda \in X^{+}$are integrable.

Definition An $L(A)$-module $V$ is called integrable if

$$
V=\bigoplus_{\lambda \in H^{*}} V_{\lambda}
$$

and if $e_{i}: V \rightarrow V$ and $f_{i}: V \rightarrow V$ are locally nilpotent for all $i=1, \ldots, n$.
Proposition 19.12 The adjoint module $L(A)$ is integrable.
Proof. The proof of Proposition 7.17 carries over to the present situation.

Proposition 19.13 Let $V$ be an integrable $L(A)$-module. Then $\operatorname{dim} V_{\lambda}=$ $\operatorname{dim} V_{w(\lambda)}$ for each $\lambda \in H^{*}$ and each $w \in W$.

Proof. Since the Weyl group $W$ of $L(A)$ is generated by the elements $s_{i}$ it is sufficient to show that $\operatorname{dim} V_{\lambda}=\operatorname{dim} V_{s_{i}(\lambda)}$.

We may regard $V$ as a module for the 3-dimensional simple subalgebra $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$ of $L(A)$. Let $v \in V_{\lambda}$ and consider the $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$-submodule generated by $v$. The vectors

$$
v, \quad e_{i} v, \quad e_{i}^{2} v, \quad \ldots, \quad e_{i}^{r-1} v
$$

lie in this submodule, where $r$ is the smallest positive integer with $e_{i}^{r} v=0$. The vectors $f_{i}^{b} e_{i}^{a} v$ also lie in this submodule, and there are only finitely many $(a, b)$ for which such a vector is non-zero. Each such vector is a weight vector in $V$. However, the relation

$$
e_{i} f_{i}^{n}=f_{i}^{n} e_{i}+n f_{i}^{n-1}\left(h_{i}-(n-1)\right)
$$

obtained in the proof of Theorem 10.20 shows that the subspace spanned by all vectors $f_{i}^{b} e_{i}^{a} v$ is an $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$-submodule. Hence every weight vector $v \in V$ lies in a finite dimensional $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$-submodule which is also an $H$-module, i.e. it is an $\left\langle e_{i}, H, f_{i}\right\rangle$-submodule.

Now let $U$ be the subspace of $V$ given by

$$
U=\sum_{k \in \mathbb{Z}} V_{\lambda+k \alpha_{i}}
$$

$U$ is clearly an $\left\langle e_{i}, H, f_{i}\right\rangle$-submodule of $V$. The $\left\langle e_{i}, H, f_{i}\right\rangle$-submodule generated by each weight vector is finite dimensional, thus $U$ is a sum of finite dimensional $\left\langle e_{i}, H, f_{i}\right\rangle$-submodules. Now $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$ is a 3-dimensional simple Lie algebra of type $A_{1}$. Thus every finite dimensional $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$-module is a direct sum of finite dimensional irreducible $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$-modules, by the complete reducibility theorem, Theorem 12.20 . The weight spaces involved in such a decomposition of an $H$-invariant $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$-module can be chosen as weight spaces for $H$, as in the proof of Theorem 10.20, thus every finite dimensional $H$-invariant $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$-module is a direct sum of finite dimensional $H$-invariant irreducible $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$-modules. Thus $U$ is a sum of finite dimensional $H$-invariant irreducible $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$-submodules, so is a direct sum of certain of these submodules. However, for each of these irreducible submodules $M$ we have

$$
\operatorname{dim} M_{\lambda}=\operatorname{dim} M_{s_{i}(\lambda)}
$$

by Proposition 10.22. It follows that

$$
\operatorname{dim} V_{\lambda}=\operatorname{dim} V_{s_{i}(\lambda)}
$$

and the required result follows.

Proposition 19.14 Let $L(A)$ be a symmetrisable Kac-Moody algebra and $L(\lambda)$ be an irreducible $L(A)$-module in the category $\mathcal{O}$. Then $L(\lambda)$ is integrable if and only if $\lambda$ is dominant and integral.

Proof. Suppose first that $L(\lambda)$ is integrable. Let $v_{\lambda}$ be a highest weight vector in $L(\lambda)$. Then $f_{i}^{r} v_{\lambda}=0$ for some $r$. Consider the vectors

$$
v_{\lambda}, \quad f_{i} v_{\lambda}, \quad \ldots, \quad f_{i}^{r-1} v_{\lambda}
$$

Since $e_{i} f_{i}^{n}=f_{i}^{n} e_{i}+n f_{i}^{n-1}\left(h_{i}-(n-1)\right)$ for each $n$ we see that these vectors span an $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$-submodule of $L(\lambda)$. The highest weight of this finite dimensional $\left\langle e_{i}, h_{i}, f_{i}\right\rangle$-module is $\lambda$. But the highest weight of any finite dimensional module for a finite dimensional simple Lie algebra is dominant and integral. Thus $\lambda\left(h_{i}\right) \in \mathbb{Z}$ and $\lambda\left(h_{i}\right) \geq 0$. Since this holds for all $i, \lambda$ is dominant and integral.

Now suppose conversely that $\lambda\left(h_{i}\right) \in \mathbb{Z}$ and $\lambda\left(h_{i}\right) \geq 0$ for each $i$. Then we have

$$
f_{i}^{\lambda\left(h_{i}\right)+1} v_{\lambda}=0
$$

as in the proof of Theorem 10.20. Now each element of $L(\lambda)$ has form $u v_{\lambda}$ for some $u \in L(A)$. We have

$$
f_{i}^{N}\left(u v_{\lambda}\right)=\sum_{k=0}^{N}\binom{N}{k}\left(\left(\operatorname{ad} f_{i}\right)^{k} u\right)\left(f_{i}^{N-k} v_{\lambda}\right)
$$

Now $\left(\operatorname{ad} f_{i}\right)^{k} u=0$ for $k$ sufficiently large since $L(A)$ is integrable, by Proposition 19.12. Also $f_{i}^{N-k} v_{\lambda}=0$ for $N-k$ sufficiently large, as shown above. Thus $f_{i}^{N}\left(u v_{\lambda}\right)=0$ for $N$ sufficiently large, and so $f_{i}: L(\lambda) \rightarrow L(\lambda)$ is locally nilpotent. The fact that $e_{i}: L(\lambda) \rightarrow L(\lambda)$ is locally nilpotent follows from the fact that $L(\lambda)$ lies in category $\mathcal{O}$. Thus $L(\lambda)$ is integrable.

As before we write

$$
X^{+}=\left\{\lambda \in H^{*} ; \lambda\left(h_{i}\right) \in \mathbb{Z}, \lambda\left(h_{i}\right) \geq 0 \quad \text { for each } i\right\}
$$

We now turn to Kac' character formula for ch $L(\lambda)$ when $\lambda \in X^{+}$. We recall from Proposition 19.2 that the character of the corresponding Verma module $M(\lambda)$ is given by

$$
\operatorname{ch} M(\lambda)=\frac{e_{\lambda}}{\Delta}
$$

where $\Delta=\prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)^{m_{\alpha}}$ and $m_{\alpha}$ is the multiplicity of $\alpha$.
We begin with a lemma.

Lemma 19.15 Let $X^{++}=\left\{\lambda \in X ; \lambda\left(h_{i}\right)>0\right.$ for all $\left.i\right\}$. Suppose $\xi \in X^{++}, \eta \in$ $X^{+}$satisfy $\eta \prec \xi$ and $\langle\eta, \eta\rangle=\langle\xi, \xi\rangle$. Then $\eta=\xi$.

Proof. Since $\eta \prec \xi$ we have $\xi-\eta=\sum_{i=1}^{n} k_{i} \alpha_{i}$ with $k_{i} \in \mathbb{Z}$ and $k_{i} \geq 0$. Thus

$$
\begin{aligned}
\langle\xi, \xi\rangle & -\langle\eta, \eta\rangle=\langle\xi+\eta, \xi-\eta\rangle \\
& =\sum_{i} k_{i}\left\langle\xi+\eta, \alpha_{i}\right\rangle=\sum_{i} k_{i} \frac{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}{2}(\xi+\eta)\left(h_{i}\right) .
\end{aligned}
$$

Now $\left\langle\alpha_{i}, \alpha_{i}\right\rangle>0$ and $(\xi+\eta)\left(h_{i}\right)>0$. Hence $\langle\xi, \xi\rangle-\langle\eta, \eta\rangle=0$ implies that $k_{i}=0$ for each $i$. Thus $\xi=\eta$.

Theorem 19.16 (Kac' character formula). Let $L(A)$ be a symmetrisable Kac-Moody algebra and $L(\lambda)$ be an irreducible $L(A)$-module with $\lambda \in X^{+}$. Then

$$
\operatorname{ch} L(\lambda)=\frac{\sum_{w \in W} \varepsilon(w) e_{w(\lambda+\rho)-\rho}}{\prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)^{m_{\alpha}}}
$$

(This is an equality in the ring $\mathfrak{R}$.

Proof. By Proposition 19.6 we have

$$
\operatorname{ch} M(\lambda)=\sum_{\mu \in H^{*}}[M(\lambda): L(\mu)] \operatorname{ch} L(\mu)
$$

Now all $\mu$ for which $[M(\lambda): L(\mu)] \neq 0$ satisfy $\mu \prec \lambda$. For $L(\mu)$ appears as a factor in some filtration of $M(\lambda)$, so $\mu$ is a weight of $M(\lambda)$.

We consider the action of the generalised Casimir operator $\Omega$ on $M(\lambda)$. By Proposition $19.10 \Omega$ acts on $M(\lambda)$ as scalar multiplication by $\langle\lambda+\rho, \lambda+\rho\rangle-$ $\langle\rho, \rho\rangle$. Similarly by Corollary $19.11 \Omega$ acts on $L(\mu)$ as scalar multiplication by $\langle\mu+\rho, \mu+\rho\rangle-\langle\rho, \rho\rangle$. Thus if $[M(\lambda): L(\mu)] \neq 0$ we must have

$$
\langle\lambda+\rho, \lambda+\rho\rangle-\langle\rho, \rho\rangle=\langle\mu+\rho, \mu+\rho\rangle-\langle\rho, \rho\rangle
$$

that is $\langle\lambda+\rho, \lambda+\rho\rangle=\langle\mu+\rho, \mu+\rho\rangle$. Thus

$$
\operatorname{ch} M(\lambda)=\sum_{\mu}[M(\lambda): L(\mu)] \operatorname{ch} L(\mu)
$$

summed over all $\mu \prec \lambda$ with $\langle\mu+\rho, \mu+\rho\rangle=\langle\lambda+\rho, \lambda+\rho\rangle$. If we take a total ordering on the weights $\mu$ satisfying $\mu \prec \lambda$ and $\langle\mu+\rho, \mu+\rho\rangle=\langle\lambda+\rho, \lambda+\rho\rangle$ which is compatible with the partial ordering $\prec$ these equations can be written

$$
\operatorname{ch} M(\lambda)=\sum_{\mu} a_{\lambda \mu} \operatorname{ch} L(\mu)
$$

where $\left(a_{\lambda \mu}\right)$ is an infinite matrix with non-negative integer entries such that $a_{\lambda \lambda}=1$ and $a_{\lambda \mu}=0$ for all entries below the diagonal. Such a matrix $\left(a_{\lambda \mu}\right)$ can be inverted to give a matrix $\left(b_{\lambda \mu}\right)$ with $b_{\lambda \mu} \in \mathbb{Z}, b_{\lambda \lambda}=1$ and $b_{\lambda \mu}=0$ for entries below the diagonal. Thus we have

$$
\begin{aligned}
\operatorname{ch} L(\lambda) & =\sum_{\mu} b_{\lambda \mu} \operatorname{ch} M(\mu) \\
& =\sum_{\mu} b_{\lambda \mu} \frac{e_{\mu}}{\Delta}
\end{aligned}
$$

Thus $\Delta \operatorname{ch} L(\lambda)=\sum_{\mu} b_{\lambda \mu} e_{\mu}$ and $e_{\rho} \Delta \operatorname{ch} L(\lambda)=\sum_{\mu} b_{\lambda \mu} e_{\mu+\rho}$.
We consider the action of the Weyl group on the functions which appear here. Since $s_{i}$ transforms $\alpha_{i}$ to $-\alpha_{i}$ and $\Phi^{+}-\left\{\alpha_{i}\right\}$ into itself we have

$$
\begin{aligned}
s_{i}\left(e_{\rho} \Delta\right) & =s_{i}\left(e_{\rho}\left(1-e_{-\alpha_{i}}\right) \prod_{\alpha \in \Phi^{+}-\left\{\alpha_{i}\right\}}\left(1-e_{-\alpha}\right)^{m_{\alpha}}\right) \\
& =e_{\rho-\alpha_{i}}\left(1-e_{\alpha_{i}}\right) \prod_{\alpha \in \Phi^{+}-\left\{\alpha_{i}\right\}}\left(1-e_{-\alpha}\right)^{m_{\alpha}} \\
& =-e_{\rho} \Delta
\end{aligned}
$$

since $s_{i}(\rho)=\rho-\alpha_{i}$. Hence

$$
w\left(e_{\rho} \Delta\right)=\varepsilon(w) e_{\rho} \Delta \quad \text { for all } w \in W
$$

Also by Proposition 19.13 we have

$$
w(\operatorname{ch} L(\lambda))=\operatorname{ch} L(\lambda) \quad \text { for all } w \in W
$$

since $L(\lambda)$ is integrable. It follows that

$$
w\left(\sum_{\mu} b_{\lambda \mu} e_{\mu+\rho}\right)=\varepsilon(w) \sum_{\mu} b_{\lambda \mu} e_{\mu+\rho}
$$

This implies that

$$
b_{\lambda \mu}=\varepsilon(w) b_{\lambda \nu}
$$

where $w(\mu+\rho)=\nu+\rho$.
Suppose $\mu$ is a weight for which $b_{\lambda \mu} \neq 0$. Consider the set of all weights $\nu$ for which $w(\mu+\rho)=\nu+\rho$ for some $w \in W$. All such weights satisfy $b_{\lambda \nu} \neq 0$ and we have $\nu \prec \lambda$. Among all such weights $\nu$ we can choose one for which the height of $\lambda-\nu$ is minimal. Then $\nu+\rho$ must lie in $X^{+}$. For if there existed an $i$ for which $(\nu+\rho)\left(h_{i}\right)<0$ we would have

$$
s_{i} w(\mu+\rho)=s_{i}(\nu+\rho)=\nu+\rho-(\nu+\rho)\left(h_{i}\right) \alpha_{i}
$$

contradicting the minimality of $\mathrm{ht}(\lambda-\nu)$. Hence $\nu+\rho \in X^{+}$. We also have

$$
\langle\nu+\rho, \nu+\rho\rangle=\langle\mu+\rho, \mu+\rho\rangle=\langle\lambda+\rho, \lambda+\rho\rangle .
$$

Thus we have

$$
\lambda+\rho \in X^{++}, \quad \nu+\rho \in X^{+}, \quad \nu+\rho \prec \lambda+\rho
$$

and $\langle\nu+\rho, \nu+\rho\rangle=\langle\lambda+\rho, \lambda+\rho\rangle$. By Lemma 19.15 this implies $\nu=\lambda$. Hence every weight $\mu$ for which $b_{\lambda \mu} \neq 0$ satisfies $\mu+\rho=w(\lambda+\rho)$ for some $w \in W$. But then

$$
b_{\lambda \mu}=\varepsilon(w) b_{\lambda \lambda}=\varepsilon(w) .
$$

Hence

$$
e_{\rho} \Delta \operatorname{ch} L(\lambda)=\sum_{w \in W} \varepsilon(w) e_{w(\lambda+\rho)} .
$$

If follows that

$$
\operatorname{ch} L(\lambda)=\frac{\sum_{w \in W} \varepsilon(w) e_{w(\lambda+\rho)-\rho}}{\Delta} .
$$

(We note that $\frac{1}{\Delta}=e_{-\lambda} \operatorname{ch} M(\lambda)$ lies in $\mathfrak{R}$.)
Corollary 19.17 (Kac' denominator formula). For a symmetrisable KacMoody algebra we have

$$
e_{\rho} \prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)^{m_{\alpha}}=\sum_{w \in W} \varepsilon(w) e_{w(\rho)} .
$$

Proof. $L(0)$ is the 1-dimensional trivial module with ch $L(0)=e_{0}$. Hence

$$
\Delta=\Delta e_{0}=\sum_{w \in W} \varepsilon(w) e_{w(\rho)-\rho}
$$

and so $e_{\rho} \Delta=\sum_{w \in W} \varepsilon(w) e_{w(\rho)}$.

Corollary 19.18 (Alternative form of Kac' character formula). Let $L(\lambda), \lambda \in X^{+}$, be an irreducible module for a symmetrisable Kac-Moody algebra. Then

$$
\operatorname{ch} L(\lambda)=\frac{\sum_{w \in W} \varepsilon(w) e_{w(\lambda+\rho)}}{\sum_{w \in W} \varepsilon(w) e_{w(\rho)}}
$$

Proof. This follows from Theorem 19.16 and Corollary 19.17.
Note Kac' character formula and denominator formula appear very similar to Weyl's character and denominator formulae for finite dimensional semisimple Lie algebras. However, the nature of Kac' formulae is in fact rather different, since they involve in general infinite sums over the elements of $W$ and infinite products over the positive roots.

Theorem 19.19 Let $L(A)$ be a symmetrisable Kac-Moody algebra and $\lambda \in X^{+}$. Then $L(\lambda)=M(\lambda) / J(\lambda)$ where $J(\lambda)$ is the submodule of $M(\lambda)$ generated by elements $f_{i}^{\lambda\left(h_{i}\right)+1} m_{\lambda} \quad$ for $i=1, \ldots, n$.

Proof. Let $K(\lambda)$ be the submodule of $M(\lambda)$ generated by the elements $f_{i}^{\lambda\left(h_{i}\right)+1} m_{\lambda}$. We know that $L(\lambda)=M(\lambda) / J(\lambda)$ where $J(\lambda)$ is the unique maximal submodule of $M(\lambda)$, and wish to show that $K(\lambda)=J(\lambda)$. Now we have

$$
f_{i}^{\lambda\left(h_{i}\right)+1} v_{\lambda}=0 \quad \text { where } v_{\lambda}=J(\lambda)+m_{\lambda}
$$

as in the proof of Proposition 19.14 (the detailed argument is given in Theorem 10.20). Thus $f_{i}^{\lambda\left(h_{i}\right)+1} m_{\lambda} \in J(\lambda)$ and so $K(\lambda) \subset J(\lambda)$.

Let $V(\lambda)=M(\lambda) / K(\lambda)$. Then $V(\lambda)$ is an $L(A)$-module in the category $\mathcal{O}$, so

$$
\operatorname{ch} V(\lambda)=\sum_{\mu<\lambda}[V(\lambda): L(\mu)] \operatorname{ch} L(\mu)
$$

by Proposition 19.6. We also have

$$
\operatorname{ch} L(\mu)=\sum_{\nu<\mu} b_{\mu \nu} \operatorname{ch} M(\nu)
$$

Hence

$$
\operatorname{ch} V(\lambda)=\sum_{\mu<\lambda} c_{\lambda \mu} \operatorname{ch} M(\mu)
$$

for certain $c_{\lambda \mu} \in \mathbb{Z}$. By considering the action of the generalised Casimir operator $\Omega$ on $M(\mu)$ and on $V(\lambda)$ and using Proposition 19.10 we have

$$
\operatorname{ch} V(\lambda)=\sum_{\substack{\mu<\lambda \\\langle\mu+\rho, \mu+\rho\rangle=\langle\lambda+\rho, \lambda+\rho\rangle}} c_{\lambda \mu} \operatorname{ch} M(\mu)
$$

Now let $v_{\lambda}^{\prime}=K(\lambda)+m_{\lambda}$ be the highest weight vector of $V(\lambda)$. Then we have

$$
f_{i}^{\lambda\left(h_{i}\right)+1} v_{\lambda}=0
$$

It follows, as in the proof of Proposition 19.14, that $f_{i}: V(\lambda) \rightarrow V(\lambda)$ is locally nilpotent. Since $V(\lambda) \in \mathcal{O}, \quad e_{i}: V(\lambda) \rightarrow V(\lambda)$ is locally nilpotent. Hence $V(\lambda)$ is an integrable $L(A)$-module. Thus

$$
w(\operatorname{ch} V(\lambda))=\operatorname{ch} V(\lambda) \quad \text { for all } w \in W
$$

by Proposition 19.13. We then have

$$
e_{\rho} \Delta \operatorname{ch} V(\lambda)=\sum_{\substack{\mu<\lambda \\\langle\mu+\rho, \mu+\rho\rangle=\langle\lambda+\rho, \lambda+\rho\rangle}} c_{\lambda \mu} e_{\mu+\rho}
$$

by Proposition 19.2. It then follows exactly as in the proof of Theorem 19.16 that every weight $\mu$ for which $c_{\lambda \mu} \neq 0$ satisfies $\mu+\rho=w(\lambda+\rho)$ for some $w \in W$, and that then $c_{\lambda \mu}=\varepsilon(w)$. Hence

$$
e_{\rho} \Delta \operatorname{ch} V(\lambda)=\sum_{w \in W} \varepsilon(w) e_{w(\lambda+\rho)}
$$

and so ch $V(\lambda)=\operatorname{ch} L(\lambda)$ by Theorem 19.16. Since $L(\lambda)$ is a factor module of $V(\lambda)$ this can only happen if $V(\lambda)=L(\lambda)$. Thus $K(\lambda)=J(\lambda)$ as required.

We now recall that for finite dimensional semisimple Lie algebras the partition function $\mathfrak{R}$ was defined as follows. If $\xi \in H^{*} \mathfrak{B}(\xi)$ is the number of ways of writing $\xi$ as a sum of positive roots, i.e. as the number of sets of non-negative integers $r_{\alpha}, \alpha \in \Phi^{+}$, such that $\xi=\sum_{\alpha \in \Phi^{+}} r_{\alpha} \alpha$.

For Kac-Moody algebras we define the generalised partition function $\Omega$ as follows. If $\xi \in H^{*} \mathscr{\Re}(\xi)$ is the number of ways of writing $\xi$ as a sum of positive roots, each such root $\alpha$ being taken $m_{\alpha}$ times, i.e. as the number of sets of non-negative integers $r_{\alpha, i}$ for $\alpha \in \Phi^{+}$and $1 \leq i \leq m_{\alpha}$ such that

$$
\xi=\sum_{\alpha \in \Phi^{+}} \sum_{i=1}^{m_{\alpha}} r_{\alpha, i} \alpha
$$

We then have an analogue for symmetrisable Kac-Moody algebras of Kostant's multiplicity formula Theorem 12.18.

Proposition 19.20 Let $L(A)$ be a symmetrisable Kac-Moody algebra and let $\lambda \in \chi^{+}$. Then for each weight $\mu$ of $L(\lambda)$ we have

$$
\operatorname{dim} L(\lambda)_{\mu}=\sum_{w \in W} \varepsilon(w) \mathscr{R}(w(\lambda+\rho)-(\mu+\rho)) .
$$

Proof. By Proposition 19.2 we have

$$
\operatorname{ch} M(\lambda)=e_{\lambda} \prod_{\alpha \in \Phi^{+}}\left(1+e_{-\alpha}+e_{-2 \alpha}+\cdots\right)^{m_{\alpha}} .
$$

By definition of $\Omega$ we have

$$
\prod_{\alpha \in \Phi^{+}}\left(1+e_{\alpha}+e_{2 \alpha}+\cdots\right)^{m_{\alpha}}=\sum_{\beta \in Q^{+}} \Omega(\beta) e_{\beta} .
$$

Thus ch $M(\lambda)=e_{\lambda} \sum_{\beta \in Q^{+}} \Re(\beta) e_{-\beta}$. It follows that

$$
\begin{aligned}
\operatorname{ch} L(\lambda) & =\sum_{w \in W} \varepsilon(w) \operatorname{ch} M(w(\lambda+\rho)-\rho) \\
& =\sum_{w \in W} \sum_{\beta \in Q^{+}} \varepsilon(w) e_{w(\lambda+\rho)-\rho} \Omega(\beta) e_{-\beta} \\
& =\sum_{w \in W} \sum_{\beta \in Q^{+}} \varepsilon(w) e_{w(\lambda+\rho)-\rho-\beta} \Omega(\beta) \\
& =\sum_{\mu} \sum_{w \in W} \varepsilon(w) \Re(w(\lambda+\rho)-(\mu+\rho)) e_{\mu} .
\end{aligned}
$$

Hence the multiplicity of $\mu$ as a weight of $L(\lambda)$ is

$$
\sum_{w \in W} \varepsilon(w) \Re(w(\lambda+\rho)-(\mu+\rho)) .
$$

### 19.4 Generators and relations for symmetrisable algebras

We recall that the Kac-Moody algebra $L(A)$ was not defined in terms of generators and relations. The larger algebra $\tilde{L}(A)$ was defined by generators and relations and its quotient $L(A)$ is given as $\tilde{L}(A) / I$ where $I$ is the largest ideal of $\tilde{L}(A)$ satisfying $I \cap \tilde{H}=O$. It is natural to ask what additional relations are required to pass from $\tilde{L}(A)$ to $L(A)$. We shall answer this in the case when the GCM $A$ is symmetrisable.

We first require some preliminary results on enveloping algebras and modules in category $\mathcal{O}$.

Proposition 19.21 Let $\theta: L \rightarrow L^{\prime}$ be a surjective homomorphism of Lie algebras with kernel $K$. Let $\phi: \mathfrak{H}(L) \rightarrow \mathfrak{U}\left(L^{\prime}\right)$ be the corresponding homomorphism between enveloping algebras. Then the kernel of $\phi$ is $K \mathfrak{H}(L)$.

Proof. Since $K$ is an ideal of $L, K \mathfrak{H}(L)$ is a 2-sided ideal of $\mathfrak{U}(L)$. For $[k x] \in$ $K$ for $k \in K, x \in L$ and so $k x=x k+[k x]$ in $\mathfrak{H}(L)$. Thus $K \mathfrak{H}(L)=\mathfrak{H}(L) K$ and $K \mathfrak{H}(L)$ is a 2-sided ideal of $\mathfrak{H}(L)$. Thus $K \mathfrak{H}(L) \subset$ ker $\phi$.

Conversely we have a homomorphism

$$
\alpha: \mathfrak{U}(L) / K \mathfrak{U}(L) \rightarrow \mathfrak{U}\left(L^{\prime}\right)
$$

induced by $\phi$. We consider the Lie algebra $[\mathfrak{H}(L) / K \mathfrak{H}(L)]$. We shall define a map $L^{\prime} \rightarrow[\mathfrak{H}(L) / K \mathfrak{H}(L)]$ as follows. Given $x^{\prime} \in L^{\prime}$ we choose $x_{1} \in L$ with $\theta\left(x_{1}\right)=x^{\prime}$. Then $x_{1} \in \mathfrak{U}(L)$ gives rise to $\bar{x}_{1} \in[\mathfrak{U}(L) / K \mathfrak{U}(L)]$. We show that the map $x^{\prime} \rightarrow \bar{x}_{1}$ is well defined. Suppose $x_{2} \in L$ also satisfies $\theta\left(x_{2}\right)=x^{\prime}$. Then $\bar{x}_{2} \in[\mathfrak{H}(L) / K \mathfrak{H}(L)]$. Now $\theta\left(x_{1}\right)=\theta\left(x_{2}\right)$ so $x_{1}-x_{2} \in K$. Hence $\bar{x}_{1}=$ $\overline{x_{1}-x_{2}}+\bar{x}_{2}=\bar{x}_{2}$. Thus our map is well defined and is clearly a Lie algebra homomorphism. By the universal property of enveloping algebras there is a homomorphism

$$
\beta: \mathfrak{U}\left(L^{\prime}\right) \rightarrow \mathfrak{U}(L) / K \mathfrak{U}(L)
$$

compatible with our homomorphism of Lie algebras

$$
L^{\prime} \rightarrow[\mathfrak{U}(L) / K \mathfrak{U}(L)] .
$$

It is readily checked that $\alpha, \beta$ are inverse homomorphisms, and thus isomorphisms. Hence the homomorphism $\phi: \mathfrak{H}(L) \rightarrow \mathfrak{H}\left(L^{\prime}\right)$ has kernel $K \mathfrak{U}(L)$.

The 2-sided ideal $L \mathfrak{H}(L)$ of $\mathfrak{H}(L)$ will be denoted by $\mathfrak{H}(L)^{+}$. We have

$$
\mathfrak{U}(L)=\mathbb{C} 1 \oplus \mathfrak{U}(L)^{+}
$$

Proposition 19.22 $L \cap\left(\mathfrak{U}(L)^{+}\right)^{2}=[L L]$.

Proof. Since $L \subset \mathfrak{H}(L)^{+}$and, for $x, y \in L,[x y]=x y-y x$ we see that

$$
[L L] \subset L \cap\left(\mathfrak{U}(L)^{+}\right)^{2}
$$

Conversely let $\bar{L}=L /[L L]$. We have a natural homomorphism $\mathfrak{H}(L) \rightarrow \mathfrak{U}(\bar{L})$ under which $L \cap\left(\mathfrak{U}(L)^{+}\right)^{2}$ maps to $\bar{L} \cap\left(\mathfrak{U}(\bar{L})^{+}\right)^{2}$. Now $\bar{L}$ is an abelian Lie
algebra so $\mathfrak{H}(\bar{L})$ is a polynomial algebra. In such a polynomial algebra it is evident that

$$
\bar{L} \cap\left(\mathfrak{U}(\bar{L})^{+}\right)^{2}=0 .
$$

It follows that $L \cap\left(\mathfrak{H}(L)^{+}\right)^{2}$ lies in the kernel of $L \rightarrow \bar{L}$ and so

$$
L \cap\left(\mathfrak{U}(L)^{+}\right)^{2} \subset[L L] .
$$

Proposition 19.23 Let $K$ be a subalgebra of the Lie algebra L. Then $K \cap K \mathfrak{H}(L)^{+}=[K K]$.

Proof. Since $K \subset \mathfrak{H}(K)^{+}$we have $[K K] \subset K \mathfrak{H}(K)^{+}$, using $[x y]=x y-y x$. Hence $[K K] \subset K \cap K \mathfrak{U}(L)^{+}$. To prove the converse we use the PBW basis theorem. Let $\left\{k_{i}\right\}$ be a basis of $K$, and extend it to a basis $\left\{k_{i}, u_{j}\right\}$ of $L$. Then all finite products of the form $\prod k_{i}^{m_{i}} u_{j}^{n_{j}}$ with $m_{i} \geq 0, n_{j} \geq 0$ form a basis of $\mathfrak{H}(L)$ and the subset $\Pi k_{i}^{m_{i}}$ with $m_{i} \geq 0$ is a basis for $\mathfrak{U}(K)$. The monomials $\Pi k_{i}^{m_{i}} u_{j}^{n_{j}}$ with $\sum m_{i}+\sum n_{j} \geq 1$ form a basis of $\mathfrak{H}(L)^{+}$and those with $\sum m_{i}+\sum n_{j} \geq 2$ and $\sum m_{i} \geq 1$ form a basis of $K \mathfrak{U}(L)^{+}$. Now

$$
K \cap K \mathfrak{U}(L)^{+} \subset \mathfrak{U}(K) \cap K \mathfrak{U}(L)^{+} .
$$

A linear combination of monomials $\prod k_{i}^{m_{i}} u_{j}^{n_{j}}$ lies in $\mathfrak{U}(K)$ if and only if all such monomials have $n_{j}=0$. Thus each element of $K \cap K \mathfrak{H}(L)^{+}$is a linear combination of such monomials with all $n_{j}=0$ and $\sum m_{i} \geq 2$. Hence

$$
K \cap K \mathfrak{U}(L)^{+} \subset\left(\mathfrak{U}(K)^{+}\right)^{2} \cap K
$$

and so $K \cap K \mathfrak{H}(L)^{+} \subset[K K]$ by Proposition 19.22.
We next need some further properties of Verma modules. We recall that for $\lambda \in H^{*}$ the Verma module $M(\lambda)$ for $L(A)$ is given by

$$
M(\lambda)=\mathfrak{l l}(L(A)) / K_{\lambda}
$$

where $K_{\lambda}=\mathfrak{l}(L(A)) N+\sum_{x \in H} \mathfrak{H}(L(A))(x-\lambda(x))$. The Verma module $M(\lambda)$ can also be described as a tensor product. Let $B$ be the subalgebra of $L(A)$ given by $B=N+H$.

Lemma 19.24 $M(\lambda)$ is isomorphic to the $\mathfrak{H}(L)$-module $\mathfrak{H}(L) \otimes_{\mathfrak{H}(B)} \mathbb{C} v_{\lambda}$, where $\mathbb{C} v_{\lambda}$ is the 1 -dimensional $B$-module with $N$ in the kernel and $H$ acting by the weight $\lambda$.

Proof. There is a bijection

$$
\mathfrak{H}(L) \otimes_{\mathfrak{l}(B)} \mathbb{C} v_{\lambda} \rightarrow \mathfrak{U}\left(N^{-}\right) \otimes_{\mathbb{C}} \mathbb{C} v_{\lambda}
$$

given as follows. Since $L=N^{-} \oplus B$ we have a bijection

$$
\mathfrak{U}(L) \rightarrow \mathfrak{U}\left(N^{-}\right) \otimes_{\mathbb{C}} \mathfrak{H}(B)
$$

Thus we have bijections

$$
\begin{aligned}
\mathfrak{H}(L) \otimes_{\mathfrak{H}(B)} \mathbb{C} v_{\lambda} & \rightarrow\left(\mathfrak{H}\left(N^{-}\right) \otimes_{\mathbb{C}} \mathfrak{H}(B)\right) \otimes_{\mathfrak{H}(B)} \mathbb{C} v_{\lambda} \\
& \rightarrow \mathfrak{H}\left(N^{-}\right) \otimes_{\mathbb{C}}\left(\mathfrak{H}(B) \otimes_{\mathfrak{H}(B)} \mathbb{C} v_{\lambda}\right) \rightarrow \mathfrak{H}\left(N^{-}\right) \otimes_{\mathbb{C}} \mathbb{C} v_{\lambda}
\end{aligned}
$$

The $\mathfrak{U}(L)$-action on $\mathfrak{U}\left(N^{-}\right) \otimes_{\mathbb{C}} \mathbb{C} v_{\lambda}$ is given as follows. Let $u^{\prime} \in \mathfrak{H}(L)$ and $u \in \mathfrak{U}\left(N^{-}\right)$. Then

$$
u^{\prime} u=\sum_{i} a_{i} b_{i} \quad \text { where } a_{i} \in \mathfrak{U}\left(N^{-}\right), b_{i} \in \mathfrak{H}(B)
$$

We have $u^{\prime}\left(u \otimes v_{\lambda}\right)=\left(\sum_{i} \lambda\left(b_{i}\right) a_{i}\right) \otimes v_{\lambda}$. On the other hand we know that each element of $M(\lambda)$ is expressible uniquely as $u m_{\lambda}$ for $u \in \mathfrak{H}\left(N^{-}\right)$. Moreover for $u^{\prime} \in \mathfrak{U}(L)$ we have

$$
u^{\prime}\left(u m_{\lambda}\right)=\left(\sum_{i} \lambda\left(b_{i}\right) a_{i}\right) m_{\lambda}
$$

Thus there is a $\mathfrak{H}(L)$-module isomorphism between $\mathfrak{H}(L) \otimes_{\mathfrak{U}(B)} \mathbb{C} v_{\lambda}$ and $M(\lambda)$.

We may also define a module $\tilde{M}(\lambda)$ for the larger Lie algebra $\tilde{L}(A)$ by

$$
\tilde{M}(\lambda)=\mathfrak{U l}(\tilde{L}) \otimes_{\mathfrak{H}(\tilde{B})} \mathbb{C} v_{\lambda}
$$

where $\lambda \in H^{*}$ and $\tilde{B}=\tilde{N}+H$.
Lemma 19.25 For $\lambda \in H^{*}$ there is an isomorphism of $\mathfrak{U}(L)$-modules

$$
\mathfrak{U}(L) \otimes_{\mathfrak{H}(\tilde{L})} \tilde{M}(\lambda) \cong M(\lambda)
$$

Proof. We have a sequence of bijections

$$
\begin{aligned}
\mathfrak{H}(L) & \otimes_{\mathfrak{H}(\tilde{L})} \tilde{M}(\lambda)=\mathfrak{H}(L) \otimes_{\mathfrak{H}(\tilde{L})}\left(\mathfrak{H}(\tilde{L}) \otimes_{\mathfrak{l}(\tilde{B})} \mathbb{C} v_{\lambda}\right) \\
& \rightarrow\left(\mathfrak{H}(L) \otimes_{\mathfrak{H}(\tilde{L})} \mathfrak{H}(\tilde{L})\right) \otimes_{\mathfrak{H}(\tilde{B})} \mathbb{C} v_{\lambda} \rightarrow \mathfrak{H}(L) \otimes_{\mathfrak{H}(\tilde{B})} \mathbb{C} v_{\lambda}
\end{aligned}
$$

Now we have a natural homomorphism $\mathfrak{U}(\tilde{B}) \rightarrow \mathfrak{U}(B)$ with kernel $K$ which acts trivially on $\mathfrak{H}(L)$ and on $\mathbb{C} v_{\lambda}$. Thus we have a bijection

$$
\mathfrak{U}(L) \otimes_{\mathfrak{H}(\tilde{B})} \mathbb{C} v_{\lambda} \rightarrow \mathfrak{H}(L) \otimes_{\mathfrak{l}(B)} \mathbb{C} v_{\lambda}=M(\lambda)
$$

The above bijections are isomorphisms of $\mathfrak{H}(L)$-modules.

We next require further information about modules in the category $\mathcal{O}$. The following definition turns out to be very useful.

Definition Let $V$ be an $L(A)$-module with $V \in \mathcal{O}$. A vector $v \in V$ is called primitive if
(i) $v$ is a weight vector
(ii) there exists a submodule $U \subset V$ such that $v \notin U$ but $N v \subset U$.

Lemma 19.26 A module $V \in \mathcal{O}$ is generated as an $L(A)$-module by its primitive vectors.

Proof. Let $V^{\prime}$ be the submodule generated by the primitive vectors in $V$. Suppose $V^{\prime} \neq V$. Consider the factor module $V / V^{\prime}$. This factor module lies in $\mathcal{O}$ so contains a weight vector $\bar{v} \neq 0$ of maximal weight with respect to $\prec$. Thus $N \bar{v}=0$. Let $v$ be a weight vector in $V$ such that $v \rightarrow \bar{v}$. Then $v \notin V^{\prime}$ and $N v \subset V^{\prime}$. Thus $v$ is a primitive vector not in $V^{\prime}$, a contradiction.

In fact the following stronger result is true.

Proposition 19.27 A module $V \in \mathcal{O}$ is generated as a $\mathfrak{U}\left(N^{-}\right)$-module by its primitive vectors.

Proof. We first show that if $v \in V$ is a weight vector which is not primitive then $v \in \mathfrak{U}\left(N^{-}\right) \mathfrak{U}(N)^{+} v$. For consider the $\mathfrak{U}(L)$-submodule of $V$ generated by $N v$. We have

$$
\begin{aligned}
\mathfrak{U}(L) N v & =\mathfrak{U}\left(N^{-}\right) \mathfrak{U}(H) \mathfrak{U}(N) N v \\
& =\mathfrak{U}\left(N^{-}\right) \mathfrak{U}(N) N v \quad \text { since } v \text { is a weight vector } \\
& =\mathfrak{U}\left(N^{-}\right) \mathfrak{U}(N)^{+} v .
\end{aligned}
$$

Now let $U$ be the $\mathfrak{U}\left(N^{-}\right)$-submodule generated by the primitive vectors in $V$. We wish to show $U=V$. We shall assume $U \neq V$ and obtain a contradiction. For each weight vector $v \in V$ we have

$$
\begin{aligned}
\mathfrak{U}(L) v & =\mathfrak{U}\left(N^{-}\right) \mathfrak{H}(H) \mathfrak{H}(N) v=\mathfrak{U}\left(N^{-}\right) \mathfrak{U}(N) v \\
& =\mathfrak{H}\left(N^{-}\right)\left(\mathbb{C} 1+\mathfrak{U}(N)^{+}\right) v \\
& =\mathfrak{H}\left(N^{-}\right) v+\mathfrak{H}\left(N^{-}\right) \mathfrak{H}(N)^{+} v .
\end{aligned}
$$

We can deduce from this that $V$ is generated as $\mathfrak{U}(L)$-module by $U$ and the $\mathfrak{U}\left(N^{-}\right)$-submodule generated by $\mathfrak{H}(N)^{+} v$ for all primitive $v \in V$. Since
$U \neq V$ there exists a primitive $v$ such that $\mathfrak{H}(N)^{+} v \not \subset U$. Let $v$ have weight $\lambda$. Then there exists a weight vector $u_{1} \in \mathfrak{H}(N)^{+}$with $u_{1} v \notin U$. So $u_{1} v$ is not primitive in $V$. Thus we have

$$
u_{1} v \in \mathfrak{U}\left(N^{-}\right) \mathfrak{U}(N)^{+} u_{1} v
$$

Hence $\mathfrak{U}(N)^{+} u_{1} v \not \subset U$. So there exists a weight vector $u_{2} \in \mathfrak{H}(N)^{+}$with $u_{2} u_{1} v \notin U$.

Continuing in this way we obtain a sequence of weight vectors $u_{1}, u_{2}, u_{3}, \ldots$ in $\mathfrak{H}(N)^{+}$such that $u_{k} \ldots u_{1} v \notin U$ for each $k$. Let the weight of $u_{i}$ be $\mu_{i}$. Then the weight of $u_{k} \ldots u_{1} v$ is $\lambda+\mu_{1}+\cdots+\mu_{k}$. We have

$$
\lambda \prec \lambda+\mu_{1} \prec \lambda+\mu_{1}+\mu_{2} \prec \cdots
$$

But $V \in \mathcal{O}$ and so such a sequence of weights must terminate after finitely many steps. This gives the required contradiction.

We now consider the module $\bigoplus_{i=1}^{n} M\left(-\alpha_{i}\right)$ in $\mathcal{O}$.
Proposition 19.28 Every primitive vector in the module $\bigoplus_{i=1}^{n} M\left(-\alpha_{i}\right)$ has weight $-\alpha$ where $\langle\alpha, \alpha\rangle=2\langle\rho, \alpha\rangle$.

Proof. Let $v$ be a primitive vector in $\bigoplus_{i=1}^{n} M\left(-\alpha_{i}\right)$ of weight $-\alpha$. Then there is a submodule $U$ of $\bigoplus_{i=1}^{n} M\left(-\alpha_{i}\right)$ such that $v \notin U$ and $N v \subset U$. Write $v=v_{1}+\cdots+v_{n}$ where $v_{i} \in M\left(-\alpha_{i}\right)$ and let $v \rightarrow \bar{v}$ where $\bar{v} \in\left(\bigoplus M\left(-\alpha_{i}\right)\right) / U$. We consider the action of the generalised Casimir operator $\Omega$ on the module $\bigoplus M\left(-\alpha_{i}\right)$. By Proposition 19.10

$$
\begin{aligned}
\Omega v_{i} & =\left(\left\langle-\alpha_{i}+\rho,-\alpha_{i}+\rho\right\rangle-\langle\rho, \rho\rangle\right) v_{i} \\
& =\left(\left\langle\alpha_{i}, \alpha_{i}\right\rangle-2\left\langle\rho, \alpha_{i}\right\rangle\right) v_{i} .
\end{aligned}
$$

Now $\rho\left(h_{i}\right)=1$ so $\left\langle\rho, \frac{2 \alpha_{i}}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle}\right\rangle=1$. Hence $\left\langle\alpha_{i}, \alpha_{i}\right\rangle=2\left\langle\rho, \alpha_{i}\right\rangle$ and so $\Omega v_{i}=0$. Thus $\Omega$ acts as 0 on $\bigoplus M\left(-\alpha_{i}\right)$. Hence $\Omega$ acts as 0 on $\left(\bigoplus M\left(-\alpha_{i}\right)\right) / U$ and $\Omega \bar{v}=0$. But $\bar{v}$ has weight $-\alpha$ and so

$$
\begin{aligned}
\Omega \bar{v} & =(\langle-\alpha+\rho,-\alpha+\rho\rangle-\langle\rho, \rho\rangle) \bar{v} \\
& =(\langle\alpha, \alpha\rangle-2\langle\rho, \alpha\rangle) \bar{v}
\end{aligned}
$$

Thus $\langle\alpha, \alpha\rangle=2\langle\rho, \alpha\rangle$ as required.
We shall now start to see the relevance of the preliminary results which we have obtained. We concentrate on the kernel $I$ of the natural homomorphism $\tilde{L}(A) \rightarrow L(A)$. We recall that $I=I^{-} \oplus I^{+}$where $I^{-} \subset \tilde{N}^{-}$and
$I^{+} \subset \tilde{N}$. We have $I^{-}=\bigoplus_{\alpha \in Q^{-}} I_{\alpha}^{-}$by Lemma 14.12. Since $\operatorname{dim} \tilde{L}(A)_{-\alpha_{i}}=$ $\operatorname{dim} L(A)_{-\alpha_{i}}=1$ we have $I_{-\alpha_{i}}^{-}=0$. Thus $I^{-}=\bigoplus_{\alpha \in Q^{-}, \alpha \neq-\alpha_{i}} I_{\alpha}^{-}$. Hence each element of $I^{-}$has form $\sum_{i=1}^{n} u_{i} f_{i}$ where $u_{i} \in \mathfrak{H}\left(\tilde{N}^{-}\right)^{+}$. It is in fact uniquely expressible in this form since $\tilde{N}^{-}$is the free Lie algebra on $f_{1}, \ldots, f_{n}$ by Proposition 14.8 and so $\mathfrak{l}\left(\tilde{N}^{-}\right)$is the free associative algebra on $f_{1}, \ldots, f_{n}$ by Proposition 9.10.

Proposition 19.29 Let $\theta: \tilde{L}(A) \rightarrow L(A)$ be the natural homomorphism with kernel $I=I^{-} \oplus I^{+}$. Then there is a homomorphism of $\tilde{L}$-modules

$$
I^{-} \rightarrow \bigoplus_{i=1}^{n} M\left(-\alpha_{i}\right)
$$

given by $\sum_{i=1}^{n} u_{i} f_{i} \rightarrow \sum_{i=1}^{n} \theta\left(u_{i}\right) m_{-\alpha_{i}}$. The kernel of this homomorphism is $\left[I^{-} I^{-}\right]$.

Proof. We begin with the $\tilde{L}$-module

$$
\tilde{M}(\lambda)=\mathfrak{H}(\tilde{L}) \underset{\mathfrak{H}(\tilde{B})}{\otimes} \mathbb{C} v_{\lambda} \quad \text { for } \lambda \in H^{*}
$$

The module $\tilde{M}(\lambda)$ has highest weight vector $\tilde{m}_{\lambda}=1 \otimes v_{\lambda}$ and, just as for the Verma module $M(\lambda)$, each element of $\tilde{M}(\lambda)$ is uniquely expressible in the form $u \tilde{m}_{\lambda}$ for $u \in \mathfrak{H}\left(\tilde{N}^{-}\right)$. We take the special case $\lambda=0$. Then $u \rightarrow u \tilde{m}_{0}$ is a bijection between $\mathfrak{H}\left(\tilde{N}^{-}\right)$and $\tilde{M}(0)$.

Now $\mathfrak{l}\left(\tilde{N}^{-}\right)$is freely generated by $f_{1}, \ldots, f_{n}$ so

$$
\mathfrak{H}\left(\tilde{N}^{-}\right)=\mathbb{C} 1 \oplus \mathfrak{H}\left(\tilde{N}^{-}\right) f_{1} \oplus \cdots \oplus \mathfrak{U}\left(\tilde{N}^{-}\right) f_{n}
$$

Thus $\bigoplus_{i=1}^{n} \mathfrak{l}\left(\tilde{N}^{-}\right) f_{i}$ is a $\mathfrak{U}\left(\tilde{N}^{-}\right)$- submodule of codimension 1 in $\mathfrak{H}\left(\tilde{N}^{-}\right)$. It corresponds to the subspace $\bigoplus_{i=1}^{n} \mathfrak{l}\left(\tilde{N}^{-}\right) f_{i} \tilde{m}_{0}$ of codimension 1 in $\tilde{M}(0)$. Let

$$
\tilde{J}(0)=\bigoplus_{i=1}^{n} \mathfrak{H}\left(\tilde{N}^{-}\right) f_{i} \tilde{m}_{0} .
$$

Then $\tilde{J}(0)$ is a $\mathfrak{U}(\tilde{L})$-submodule of $\tilde{M}(0)$. For it is clearly invariant under $\mathfrak{H}\left(\tilde{N}^{-}\right)$and $\mathfrak{H}(H)$, but also

$$
\begin{aligned}
& e_{i} f_{i} \tilde{m}_{0}=f_{i} e_{i} \tilde{m}_{0}+h_{i} \tilde{m}_{0}=0 \\
& e_{j} f_{i} \tilde{m}_{0}=f_{i} e_{j} \tilde{m}_{0}=0 \quad \text { if } j \neq i .
\end{aligned}
$$

Thus $\mathfrak{U}(\tilde{N})^{+} f_{i} \tilde{m}_{0}=0$ and so $\mathfrak{U}(\tilde{N}) f_{i} \tilde{m}_{0}=\mathbb{C} f_{i} \tilde{m}_{0}$. Hence

$$
\begin{aligned}
\mathfrak{H}(\tilde{L}) f_{i} \tilde{m}_{0} & =\mathfrak{H}\left(\tilde{N}^{-}\right) \mathfrak{H}(H) \mathfrak{H}(\tilde{N}) f_{i} \tilde{m}_{0} \\
& =\mathfrak{U}\left(\tilde{N}^{-}\right) \mathfrak{U}(H) f_{i} \tilde{m}_{0}=\mathfrak{H}\left(\tilde{N}^{-}\right) f_{i} \tilde{m}_{0} .
\end{aligned}
$$

Thus $\tilde{J}(0)=\bigoplus_{i=1}^{n} \mathfrak{H}(\tilde{L}) f_{i} \tilde{m}_{0}$, which is a $\mathfrak{H}(\tilde{L})$-submodule of $\tilde{M}(0)$. Now $f_{i} \tilde{m}_{0}$ has weight $-\alpha_{i}$ and the map

$$
\begin{aligned}
& \mathfrak{U}\left(\tilde{N}^{-}\right) f_{i} \tilde{m}_{0} \rightarrow \tilde{M}\left(-\alpha_{i}\right) \\
& u_{i} f_{i} \tilde{m}_{0} \rightarrow u_{i} \tilde{m}_{-\alpha_{i}}
\end{aligned}
$$

is an isomorphism of $\mathfrak{U}(\tilde{L})$-modules. Thus $\tilde{J}(0)$ is isomorphic to $\bigoplus_{i=1}^{n} \tilde{M}\left(-\alpha_{i}\right)$ as $\mathfrak{U}(\tilde{L})$-modules. It then follows from Lemma 19.25 that

$$
\mathfrak{H}(L) \otimes_{\mathfrak{H}(\tilde{L})} \tilde{J}(0) \cong \bigoplus_{i=1}^{n} M\left(-\alpha_{i}\right)
$$

as $\mathfrak{U}(\tilde{L})$-modules, or as $\mathfrak{U}(L)$-modules.
We now consider the map

$$
\phi: I^{-} \rightarrow \mathfrak{H}(L) \underset{\mathfrak{H}(\tilde{L})}{\otimes} \tilde{J}(0)
$$

given by $x \rightarrow 1 \otimes x \tilde{m}_{0}$. We note that $x \tilde{m}_{0}$ lies in $\tilde{J}(0)$ since $I^{-} \subset \tilde{N}^{-}$. We show that $\phi$ is a homomorphism of $\tilde{L}$-modules. To see this let $y \in \tilde{L}$. Then

$$
\begin{aligned}
{[y, x] } & \rightarrow 1 \otimes[y, x] \tilde{m}_{0} \\
& =1 \otimes(y x-x y) \tilde{m}_{0} \\
& =1 \otimes y\left(x \tilde{m}_{0}\right)-1 \otimes x\left(y \tilde{m}_{0}\right)
\end{aligned}
$$

Now we have $x \tilde{m}_{0} \in \tilde{J}(0)$ and $y \tilde{m}_{0} \in \tilde{J}(0)$. Thus

$$
\begin{aligned}
{[y, x] } & \rightarrow \theta(y) \otimes\left(x \tilde{m}_{0}\right)-\theta(x) \otimes\left(y \tilde{m}_{0}\right) \\
& =\theta(y) \otimes\left(x \tilde{m}_{0}\right) \quad \text { since } \theta(x)=0 \\
& =y\left(1 \otimes x \tilde{m}_{0}\right)
\end{aligned}
$$

This shows that $\phi$ is a homomorphism of $\tilde{L}$-modules. Moreover [ $I^{-} I^{-}$] lies in the kernel of $\phi$. For if $x, y \in I^{-}$we have $[y, x] \rightarrow 0$ as above, since $\theta(x)=$ $\theta(y)=0$. Thus we have a homomorphism of $\tilde{L}$-modules

$$
\phi: I^{-} \rightarrow \bigoplus_{i=1}^{n} M\left(-\alpha_{i}\right)
$$

with $\sum_{i=1}^{n} u_{i} f_{i} \rightarrow \sum_{i=1}^{n} \theta\left(u_{i}\right) m_{-\alpha_{i}}$ where $u_{i} \in \mathfrak{H}\left(\tilde{N}^{-}\right)^{+}$.
We determine the kernel $K$ of $\phi$. We know that $\left[I^{-} I^{-}\right] \subset K$ and prove the reverse inclusion. Let $\sum u_{i} f_{i} \in K$ where $u_{i} \in \mathfrak{U}\left(\tilde{N}^{-}\right)^{+}$. Then $\sum \theta\left(u_{i}\right) m_{-\alpha_{i}}=0$. This implies $\theta\left(u_{i}\right) m_{-\alpha_{i}}=0$ for each $i$ and then that $\theta\left(u_{i}\right)=0$ for each $i$. Now the homomorphism of Lie algebras $\theta: \tilde{N}^{-} \rightarrow N^{-}$gives rise to a homomorphism of enveloping algebras $\mathfrak{H}\left(\tilde{N}^{-}\right) \rightarrow \mathfrak{U}\left(N^{-}\right)$with kernel $I^{-} \mathfrak{U}\left(\tilde{N}^{-}\right)$
by Proposition 19.21. Thus $u_{i} \in I^{-} \mathfrak{U}\left(\tilde{N}^{-}\right)$and so $\sum u_{i} f_{i} \in I^{-} \mathfrak{H}\left(\tilde{N}^{-}\right)^{+}$. Hence $K \subset I^{-} \cap I^{-} \mathfrak{U}\left(\tilde{N}^{-}\right)^{+}$. However, $I^{-} \cap I^{-} \mathfrak{U}\left(\tilde{N}^{-}\right)^{+}=\left[I^{-} I^{-}\right]$by Proposition 19.23. Hence $K \subset\left[I^{-} I^{-}\right]$. Thus the kernel of our homomorphism is $\left[I^{-} I^{-}\right]$.

We now come to our description of $L(A)$ by generators and relations.

Theorem 19.30 Let $L(A)$ be a symmetrisable Kac-Moody algebra. Then $L(A)=\tilde{L}(A) / J$ where $J$ is the ideal of $\tilde{L}(A)$ generated by the elements $\left(\operatorname{ad} e_{i}\right)^{1-A_{i j}} e_{j}$ and $\left(\operatorname{ad} f_{i}\right)^{1-A_{i j}} f_{j}$ for all $i \neq j$. Thus we obtain a system of generators and relations for $L(A)$ by taking generators and relations for $\tilde{L}(A)$ and adding the further relations

$$
\left(\operatorname{ad} e_{i}\right)^{1-A_{i j}} e_{j}=0, \quad\left(\operatorname{ad} f_{i}\right)^{1-A_{i j}} f_{j}=0
$$

for all $i \neq j$.
Proof. Let $J$ be the ideal of $\tilde{L}(A)$ generated by the elements $\left(\operatorname{ad} e_{i}\right)^{1-A_{i j}} e_{j}$ and $\left(\operatorname{ad} f_{i}\right)^{1-A_{i j}} f_{j}$. We have $L(A)=\tilde{L}(A) / I$ and $J \subset I$ by Proposition 16.10. We wish to show that $I=J$.

We shall suppose if possible that $I \neq J$ and obtain a contradiction. Let $\bar{I}=I / J$. Then $\bar{I} \neq 0$ and $\bar{I}=\bar{I}^{+} \oplus \bar{I}^{-}$where

$$
\bar{I}^{+}=\bigoplus_{\alpha \in Q^{+}} \bar{I}_{\alpha}, \quad \bar{I}^{-}=\bigoplus_{\alpha \in Q^{-}} \bar{I}_{\alpha}
$$

since the analogous property holds for $I$. The automorphism $\tilde{\omega}$ of $\tilde{L}(A)$ given in Proposition 14.5 satisfies $\tilde{\omega}(I)=I$ and $\tilde{\omega}(J)=J$ so induces an automorphism on $\bar{I}=I / J$. This automorphism satisfies $\tilde{\omega}\left(\bar{I}^{+}\right)=\bar{I}^{-}$. Hence $\bar{I}^{+}=0$ if and only if $\bar{I}^{-}=0$. Since $\bar{I}=\bar{I}^{+} \oplus \bar{I}^{-}$and $\bar{I} \neq 0$ we must have $\bar{I}^{-} \neq 0$.

We know from Section 16.2 that the Weyl group $W$ acts on the weights of $\tilde{L}(A) / I$ and that weights in the same $W$-orbit have the same multiplicity. The same argument can be applied to $\tilde{L}(A) / J$ to give a similar result. Since

$$
\operatorname{dim}(\tilde{L}(A) / J)_{\alpha}=\operatorname{dim}(\tilde{L}(A) / I)_{\alpha}+\operatorname{dim}(I / J)_{\alpha}
$$

we see that $W$ acts on the weights of $\bar{I}$ and that weights in the same $W$-orbit have the same multiplicity. In fact $W$ acts on the weights of $\bar{I}^{-}$since if $\alpha \in Q^{-}$ is a weight of $\bar{I}^{-}$then $s_{i}(\alpha) \in Q^{-}$also, since $-\alpha_{i}$ is not a weight of $I$.

We choose a weight $\alpha=\sum_{i=1}^{n} k_{i} \alpha_{i} \in Q^{+}$such that $\bar{I}_{-\alpha}^{-} \neq 0$ and $\alpha$ has minimal possible height $\sum k_{i}$. Since $\bar{I}_{-s_{i}(\alpha)}^{-} \neq 0$ we have ht $s_{i}(\alpha) \geq$ ht $\alpha$. Since

$$
s_{i}(\alpha)=\alpha-2 \frac{\left\langle\alpha_{i}, \alpha\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} \alpha_{i}
$$

we have $\left\langle\alpha_{i}, \alpha\right\rangle \leq 0$. Since $\alpha=\sum k_{i} \alpha_{i}$ with each $k_{i} \geq 0$ we deduce that $\langle\alpha, \alpha\rangle \leq 0$. On the other hand we have

$$
2\langle\rho, \alpha\rangle=\sum_{i} k_{i}\left\langle\rho, 2 \alpha_{i}\right\rangle=\sum_{i} k_{i}\left\langle\alpha_{i}, \alpha_{i}\right\rangle>0 .
$$

Thus $\langle\alpha, \alpha\rangle \leq 0$ and $2\langle\rho, \alpha\rangle>0$. In particular $\langle\alpha, \alpha\rangle \neq 2\langle\rho, \alpha\rangle$. Thus the weights $-\alpha$ for $\bar{I}^{-}$for which $\alpha$ has minimal height satisfy $\langle\alpha, \alpha\rangle \neq 2\langle\rho, \alpha\rangle$.

We now recall from Proposition 19.29 that $I^{-} /\left[I^{-}, I^{-}\right]$is isomorphic as $\tilde{L}$-module to a submodule of $\bigoplus_{i=1}^{n} M\left(-\alpha_{i}\right)$. By Proposition 19.28 all primitive vectors in $\bigoplus_{i=1}^{n} M\left(-\alpha_{i}\right)$ have a weight $-\alpha$ satisfying $\langle\alpha, \alpha\rangle=2\langle\rho, \alpha\rangle$. Thus all primitive vectors of $I^{-} /\left[I^{-} I^{-}\right]$have weight $-\alpha$ satisfying $\langle\alpha, \alpha\rangle=$ $2\langle\rho, \alpha\rangle$. Now $I^{-} /\left[I^{-} I^{-}\right]$is generated as an $N^{-}$-module by its primitive vectors, by Proposition 19.27. Thus $I^{-} /\left[I^{-} I^{-}\right]$is generated as an $\tilde{N}^{-}$-module by its weight vectors with weight $-\alpha$ satisfying $\langle\alpha, \alpha\rangle=2\langle\rho, \alpha\rangle$. (Recall that $\tilde{N}^{-} / I^{-} \cong N^{-}$.)

We claim the same is true of $I^{-}$. Let $K$ be the $\tilde{N}^{-}$-submodule of $I^{-}$generated by all weight vectors with weight $-\alpha$ satisfying $\langle\alpha, \alpha\rangle=$ $2\langle\rho, \alpha\rangle$. Then $\left(\left[I^{-} I^{-}\right]+K\right) /\left[I^{-} I^{-}\right]$has the same property in $I^{-} /\left[I^{-} I^{-}\right]$, thus $\left[I^{-} I^{-}\right]+K=I^{-}$. Suppose if possible that $K \neq I^{-}$. Then $I^{-} / K$ is an $\tilde{N}^{-}$-module whose weights are non-zero elements of $Q^{-}$. Consider the submodule $\left[I^{-} / K, I^{-} / K\right]$ of $I^{-} / K$. This is an $\tilde{N}^{-}$-module whose weights have form $\beta+\gamma$ where $\beta, \gamma$ are weights of $I^{-} / K$. Thus if $\alpha$ is a weight of $I^{-} / K$ for which $\mid$ ht $\alpha \mid$ is minimal then $\alpha$ cannot be a weight of $\left[I^{-} / K, I^{-} / K\right]$. Thus

$$
\left[I^{-} / K, I^{-} / K\right] \neq I^{-} / K
$$

and this gives $K+\left[I^{-}, I^{-}\right] \neq I^{-}$, a contradiction. Thus $I^{-}$is generated as $\tilde{N}^{-}$-module by its weight vectors with weight $-\alpha$ satisfying $\langle\alpha, \alpha\rangle=2\langle\rho, \alpha\rangle$. The same must therefore be true of $\bar{I}^{-}$. However, we have seen above that the weights $-\alpha$ of $\bar{I}^{-}$for which ht $\alpha$ is minimal do not satisfy $\langle\alpha, \alpha\rangle=$ $2\langle\rho, \alpha\rangle$. This implies that the set of weight vectors with weight $-\alpha$ satisfying $\langle\alpha, \alpha\rangle=2\langle\rho, \alpha\rangle$ cannot generate $\bar{I}^{-}$as $\tilde{N}^{-}$-module. This gives the required contradiction.

## 20

## Representations of affine Kac-Moody algebras

### 20.1 Macdonald's identities

We now consider Kac' denominator formula

$$
\prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)^{m_{\alpha}}=\sum_{w \in W} \varepsilon(w) e_{w(\rho)-\rho}
$$

in the special case when $L$ is an affine Kac-Moody algebra.
We assume first that $L$ is an untwisted affine algebra. Then $L=\hat{\mathfrak{R}}\left(L^{0}\right)$ where $L^{0}$ is a finite dimensional simple Lie algebra with root system $\Phi^{0}$ and Weyl group $W^{0}$. We recall from Theorems 17.18 and 16.27 , and Corollary 18.6 that

$$
\Phi=\left\{\alpha+n \delta ; \alpha \in \Phi^{0}, n \in \mathbb{Z}\right\} \cup\{n \delta ; n \in \mathbb{Z}, n \neq 0\}
$$

and that $\alpha+n \delta$ has multiplicity 1 and $n \delta$ has multiplicity $l$. Also

$$
\Phi^{+}=\left\{\alpha+n \delta ; \alpha \in \Phi^{0}, n>0\right\} \cup\left(\Phi^{0}\right)^{+} \cup\{n \delta ; n>0\}
$$

Thus the left-hand side of the denominator formula can be expressed as

$$
\prod_{\alpha \in\left(\Phi^{0}\right)^{+}}\left(1-e_{-\alpha}\right) \prod_{n>0}\left\{\left(1-e_{-n \delta}\right)^{l} \prod_{\alpha \in \Phi^{0}}\left(1-e_{-\alpha-n \delta}\right)\right\}
$$

We also recall from Remark 17.34 that $W=t\left(M^{*}\right) W^{0}$ where $M^{*}$ is the lattice given by

$$
M^{*}= \begin{cases}\sum_{i=1}^{l} \mathbb{Z} \alpha_{i} & \text { for types } \tilde{A}_{l}, \tilde{D}_{l}, \tilde{E}_{6}, \tilde{E}_{7}, \tilde{E}_{8} \\ \sum_{\alpha_{i} \text { long }} \mathbb{Z} \alpha_{i}+\sum_{\alpha_{i} \text { short }} p \mathbb{Z} \alpha_{i} & \text { for } \tilde{B}_{l}, \tilde{C}_{l}, \tilde{F}_{4}, \tilde{G}_{2}\end{cases}
$$

and $t\left(M^{*}\right)$ is the set of $t_{\alpha}: H^{*} \rightarrow H^{*}$ for $\alpha \in M^{*}$ given by

$$
t_{\alpha}(\lambda)=\lambda+\lambda(c) \alpha-\left(\langle\lambda, \alpha\rangle+\frac{1}{2}\langle\alpha, \alpha\rangle \lambda(c)\right) \delta
$$

In calculating the right-hand side of the denominator formula we recall that

$$
\begin{aligned}
& H=H^{0} \oplus(\mathbb{C} c+\mathbb{C} d) \\
& H^{*}=\left(H^{0}\right)^{*} \oplus(\mathbb{C} \gamma+\mathbb{C} \delta)
\end{aligned}
$$

where $\left(H^{0}\right)^{*}$ is embedded in $H^{*}$ by assuming $\lambda(c)=0, \lambda(d)=0$ for $\lambda \in\left(H^{0}\right)^{*}$.

Lemma 20.1 Let $\lambda \in H^{*}$. Then

$$
\lambda=\lambda^{0}+\lambda(c) \gamma+a_{0}^{-1} \lambda(d) \delta
$$

where $\lambda^{0} \in\left(H^{0}\right)^{*}$.

Proof. Let $\lambda=\lambda^{0}+r \gamma+s \delta$ where $\lambda^{0} \in\left(H^{0}\right)^{*}$. Then $\lambda(c)=\lambda^{0}(c)+r \gamma(c)+$ $s \delta(c)$. But we have $\alpha_{i}(c)=0$ for $i=1, \ldots, l$ hence $\lambda^{0}(c)=0$. Also $\delta(c)=0$ and $\gamma(c)=1$. Hence $r=\lambda(c)$. Also $\lambda(d)=\lambda^{0}(d)+r \gamma(d)+s \delta(d)$. We know $\alpha_{i}(d)=0$ for $i=1, \ldots, l$ thus $\lambda^{0}(d)=0$. Also $\gamma(d)=0$ and $\delta(d)=a_{0}$. Hence $\lambda(d)=a_{0} s$ and $s=a_{0}^{-1} \lambda(d)$.

Of course in the untwisted case we have $a_{0}=1$.
We recall that $\rho \in H^{*}$ satisfies $\rho\left(h_{i}\right)=1$ for $i=0,1, \ldots, l$ and $\rho(d)=0$. In particular we have $\rho(c)=c_{0}+c_{1}+\cdots+c_{l}$.

Definition The number $h=a_{0}+a_{1}+\cdots+a_{l}$ is called the Coxeter number of $L$. The number $h^{\mathrm{v}}=c_{0}+c_{1}+\cdots+c_{l}$ is called the dual Coxeter number of $L$.

We note that if $L=\hat{\mathfrak{R}}\left(L^{0}\right)$ is of untwisted type the Coxeter number of $L$ is equal to the Coxeter number of $L^{0}$.

The values of $h$ and $h^{\mathrm{v}}$ are given in the following table.

| Type of $L$ | $\underline{h}$ | $\underline{h^{v}}$ |
| :---: | :---: | :---: |
| $\tilde{A}_{l}$ | $l+1$ | $l+1$ |
| $\tilde{B}_{l}$ | $2 l$ | $2 l-1$ |
| $\tilde{C}_{l}$ | $2 l$ | $l+1$ |
| $\tilde{D}_{l}$ | $2 l-2$ | $2 l-2$ |
| $\tilde{E}_{6}$ | 12 | 12 |
| $\tilde{E}_{7}$ | 18 | 18 |
| $\tilde{E}_{8}$ | 30 | 30 |
| $\tilde{F}_{4}$ | 12 | 9 |
| $\tilde{G}_{2}$ | 6 | 4 |
| $\tilde{B}_{l}^{\text {t }}$ | $2 l-1$ | $2 l$ |
| $\tilde{C}_{l}^{\text {t }}$ | $l+1$ | $2 l$ |
| $\tilde{F}_{4}^{\text {t }}$ | 9 | 12 |
| $\tilde{G}_{2}^{\mathrm{t}}$ | 4 | 6 |
| $\tilde{A}_{1}^{\prime}$ | 3 | 3 |
| $\tilde{C}_{l}^{\prime}$ | $2 l+1$ | $2 l+1$ |

Lemma $20.2 \rho=\rho^{0}+h^{v} \gamma$ where $\rho^{0} \in\left(H^{0}\right)^{*}$ satisfies $\rho^{0}\left(h_{i}\right)=1$ for $i=$ $1, \ldots, l$.

Proof. This follows from Lemma 20.1 since $\rho(c)=h^{\mathrm{v}}$ and $\rho(d)=0$.

We now consider the right-hand side of the denominator formula. Let $w \in W$ have form $w=w^{0} t_{\alpha}$ where $w^{0} \in W^{0}$ and $\alpha \in M^{*}$. Then

$$
\begin{aligned}
w(\rho)-\rho & =w^{0} t_{\alpha}(\rho)-\rho \\
& =w^{0}\left(\rho+h^{\mathrm{v}} \alpha-\left(\langle\rho, \alpha\rangle+\frac{1}{2}\langle\alpha, \alpha\rangle h^{\mathrm{v}}\right) \delta\right)-\rho \\
& =w^{0}(\rho)-\rho+h^{\mathrm{v}} w^{0}(\alpha)-\left(\langle\rho, \alpha\rangle+\frac{1}{2}\langle\alpha, \alpha\rangle h^{\mathrm{v}}\right) \delta \\
& =w^{0}\left(\rho^{0}\right)-\rho^{0}+h^{\mathrm{v}} w^{0}(\alpha)-\left(\left\langle\rho^{0}, \alpha\right\rangle+\frac{1}{2}\langle\alpha, \alpha\rangle h^{\mathrm{v}}\right) \delta
\end{aligned}
$$

since $w^{0}(\gamma)=\gamma$ and $\langle\gamma, \alpha\rangle=0$ for all $\alpha \in M^{*}$

$$
=w^{0}\left(h^{\mathrm{v}} \alpha+\rho^{0}\right)-\rho^{0}-\frac{\left(\left\langle\rho^{0}+h^{\mathrm{v}} \alpha, \rho^{0}+h^{\mathrm{v}} \alpha\right\rangle-\left\langle\rho^{0}, \rho^{0}\right\rangle\right)}{2 h^{\mathrm{v}}} \delta
$$

For convenience we shall write, for $\lambda \in\left(H^{0}\right)^{*}$,

$$
c(\lambda)=\left\langle\lambda+\rho^{0}, \lambda+\rho^{0}\right\rangle-\left\langle\rho^{0}, \rho^{0}\right\rangle
$$

We recall from Corollary 19.11 that when $\lambda$ is dominant and integral the generalised Casimir operator acts on the irreducible module with highest weight $\lambda$ as scalar multiplication by $c(\lambda)$.

We also write, for $\lambda \in\left(H^{0}\right)^{*}$,

$$
\chi^{0}(\lambda)=\frac{\sum_{w \in W^{0}} \varepsilon(w) e_{w\left(\lambda+\rho^{0}\right)-\rho^{0}}}{\sum_{w \in W^{0}} \varepsilon(w) e_{w\left(\rho^{0}\right)-\rho^{0}}}
$$

We recall from Theorem 12.17 that when $\lambda$ is dominant and integral $\chi^{0}(\lambda)$ is the character of the irreducible $L^{0}$-module $L(\lambda)$. However, $c(\lambda)$ and $\chi^{0}(\lambda)$ are now defined for all $\lambda \in\left(H^{0}\right)^{*}$. Then we have, writing $e(\lambda)$ instead of $e_{\lambda}$ for convenience:

$$
\begin{aligned}
\sum_{w \in W} & \varepsilon(w) e(w(\rho)-\rho) \\
& =\sum_{\alpha \in M^{*}} \sum_{w^{0} \in W^{0}} \varepsilon\left(w^{0}\right) e\left(w^{0}\left(h^{\mathrm{v}} \alpha+\rho^{0}\right)-\rho^{0}\right) e\left(\frac{-c\left(h^{\mathrm{v}} \alpha\right)}{2 h^{\mathrm{v}}} \delta\right) \\
& =\sum_{w^{0} \in W^{0}} \varepsilon\left(w^{0}\right) e\left(w^{0}\left(\rho^{0}\right)-\rho^{0}\right) \sum_{\alpha \in M^{*}} \chi^{0}\left(h^{\mathrm{v}} \alpha\right) e\left(\frac{-c\left(h^{\mathrm{v}} \alpha\right)}{2 h^{\mathrm{v}}} \delta\right) \\
& =\prod_{\alpha \in\left(\Phi^{0}\right)^{+}}\left(1-e_{-\alpha}\right) \sum_{\alpha \in M^{*}} \chi^{0}\left(h^{\mathrm{v}} \alpha\right) e\left(\frac{-c\left(h^{\mathrm{v}} \alpha\right)}{2 h^{\mathrm{v}}} \delta\right)
\end{aligned}
$$

by Weyl's denominator formula.
We now put $q=e_{-\delta}$ and equate the left- and right-hand sides of Kac' denominator formula. We obtain the following result.

Theorem 20.3 (Macdonald's identity for untwisted affine Kac-Moody algebras).

$$
\prod_{n>0}\left\{\left(1-q^{n}\right)^{l} \prod_{\alpha \in \Phi^{0}}\left(1-q^{n} e_{-\alpha}\right)\right\}=\sum_{\alpha \in M^{*}} \chi^{0}\left(h^{\mathrm{v}} \alpha\right) q^{c\left(h^{v} \alpha\right) / 2 h^{v}}
$$

where

$$
M^{*}= \begin{cases}\sum_{i=1}^{l} \mathbb{Z} \alpha_{i} & \text { for types } \tilde{A}_{l}, \tilde{D}_{l}, \tilde{E}_{6}, \tilde{E}_{7}, \tilde{E}_{8} \\ \sum_{\alpha_{i} \text { long }} \mathbb{Z} \alpha_{i}+\sum_{\alpha_{i} \text { short }} p \mathbb{Z} \alpha_{i} & \text { for } \tilde{B}_{l}, \tilde{C}_{l}, \tilde{F}_{4}, \tilde{G}_{2}\end{cases}
$$

We next wish to state Macdonald's identities for the twisted affine KacMoody algebras. The left-hand side of the identity is obtained from a knowledge of the real and imaginary roots together with the multiplicities of the imaginary roots. The real roots are given in Theorem 17.8 and the multiplicities of the imaginary roots in Corollaries 18.10 and 18.15. The right-hand side of the identity looks the same as before - the only change being that the appropriate lattice $M^{*}$ must be taken in each case. The appropriate lattice was described in Remark 17.34.

Theorem 20.4 (Macdonald's identity for twisted affine Kac-Moody algebras).
(a) The left-hand side of the identity is given as follows.

$$
\begin{aligned}
\tilde{B}_{l}^{\mathrm{t}} \quad & \prod_{n>0}\left\{\left(1-q^{2 n}\right)^{l}\left(1-q^{2 n-1}\right)^{l-1} \prod_{\alpha \in \Phi_{\mathrm{s}}^{0}}\left(1-q^{n} e_{-\alpha}\right) \prod_{\alpha \in \Phi_{1}^{0}}\left(1-q^{2 n} e_{-\alpha}\right)\right\} \\
\tilde{C}_{l}^{\mathrm{t}} \quad & \prod_{n>0}\left\{\left(1-q^{2 n}\right)^{l}\left(1-q^{2 n-1}\right) \prod_{\alpha \in \Phi_{\mathrm{s}}^{0}}\left(1-q^{n} e_{-\alpha}\right) \prod_{\alpha \in \Phi_{1}^{0}}\left(1-q^{2 n} e_{-\alpha}\right)\right\} \\
\tilde{F}_{4}^{\mathrm{t}} \quad & \prod_{n>0}\left\{\left(1-q^{2 n}\right)^{4}\left(1-q^{2 n-1}\right)^{2} \prod_{\alpha \in \Phi_{\mathrm{s}}^{0}}\left(1-q^{n} e_{-\alpha}\right) \prod_{\alpha \in \Phi_{1}^{0}}\left(1-q^{2 n} e_{-\alpha}\right)\right\} \\
\tilde{G}_{2}^{\mathrm{t}} \quad & \prod_{n>0}\left\{\left(1-q^{3 n}\right)^{2}\left(1-q^{3 n-1}\right)\left(1-q^{3 n-2}\right)\right. \\
& \left.\prod_{\alpha \in \Phi_{\mathrm{s}}^{0}}\left(1-q^{n} e_{-\alpha}\right) \prod_{\alpha \in \Phi_{1}^{0}}\left(1-q^{3 n} e_{-\alpha}\right)\right\} \\
\tilde{C}_{l}^{\prime} \quad & \prod_{n>0}\left\{\left(1-q^{n}\right)^{l} \prod_{\alpha \in \Phi_{\mathrm{s}}^{0}}\left(1-q^{n} e_{-\alpha}\right)\right. \\
& \left.\prod_{\alpha \in \Phi_{1}^{0}}\left(1-q^{\frac{2 n-1}{2}} e_{\frac{-1}{2} \alpha}\right)\left(1-q^{2 n} e_{-\alpha}\right)\right\} \\
\tilde{A}_{1}^{\prime} \quad & \prod_{n>0}\left\{\left(1-q^{n}\right) \prod_{\alpha \in \Phi^{0}}\left(1-q^{\frac{2 n-1}{2}} e_{-1 / 2 \alpha}\right)\left(1-q^{2 n} e_{-\alpha}\right)\right\} .
\end{aligned}
$$

(b) The right-hand side of the identity is

$$
\sum_{\alpha \in M^{*}} \chi^{0}\left(h_{\alpha}^{\mathrm{v}}\right) q^{c\left(h^{\mathrm{v}} \alpha\right) / 2 h^{\mathrm{v}}}
$$

where

$$
M^{*}= \begin{cases}\sum_{i=1}^{l} \mathbb{Z} \alpha_{i} & \text { for types } \tilde{B}_{l}^{\mathrm{t}}, \tilde{C}_{l}^{\mathrm{t}}, \tilde{F}_{4}^{\mathrm{t}}, \tilde{G}_{2}^{\mathrm{t}} \\ \sum_{\alpha_{i} \text { long }} \frac{1}{2} \mathbb{Z} \alpha_{i}+\sum_{\alpha_{i} \text { short }} \mathbb{Z} \alpha_{i} & \text { for type } \tilde{C}_{l}^{\prime} \\ \frac{1}{2} \mathbb{Z} \alpha_{1} & \text { for type } \tilde{A}_{1}^{\prime}\end{cases}
$$

We now give some examples to illustrate Macdonald's identity. Suppose first that $L$ has type $\tilde{A}_{1}$. Then $L^{0}$ has type $A_{1}, \Phi^{0}=\left\{\alpha_{1},-\alpha_{1}\right\}, h^{v}=2$ and $M^{*}=\mathbb{Z} \alpha_{1}$. Moreover $\left\langle\alpha_{1}, \alpha_{1}\right\rangle=2 \frac{c_{1}}{a_{1}}=2$ and $\rho^{0}=\frac{1}{2} \alpha_{1}$. Let $z=e_{-\alpha_{1}}$. The lefthand side of Macdonald's identity is

$$
\prod_{n>0}\left(1-q^{n}\right)\left(1-q^{n} z\right)\left(1-q^{n} z^{-1}\right)
$$

Now $\chi^{0}\left(n \alpha_{1}\right)=\frac{e_{n \alpha_{1}}-e_{-(n+1) \alpha_{1}}}{1-e_{-\alpha_{1}}}=\frac{z^{-n}-z^{n+1}}{1-z}$. We also have

$$
\begin{aligned}
c\left(2 n \alpha_{1}\right) & =\left\langle\left(2 n+\frac{1}{2}\right) \alpha_{1},\left(2 n+\frac{1}{2}\right) \alpha_{1}\right\rangle-\left\langle\frac{1}{2} \alpha_{1}, \frac{1}{2} \alpha_{1}\right\rangle \\
& =4 n(2 n+1) .
\end{aligned}
$$

Thus the right-hand side of Macdonald's identity is

$$
\sum_{n \in \mathbb{Z}} \frac{z^{-2 n}-z^{2 n+1}}{1-z} q^{n(2 n+1)}
$$

This can be written in the convenient form

$$
\frac{1}{1-z} \sum_{n \in \mathbb{Z}}\left(z^{-2 n} q^{n(2 n+1)}-z^{2 n+1} q^{n(2 n+1)}\right)=\frac{1}{1-z} \sum_{m \in \mathbb{Z}}(-1)^{m} z^{m} q^{m(m-1) / 2}
$$

Multiplying both sides of the identity by $1-z$ we obtain:

Proposition 20.5 (Macdonald's identity for type $\tilde{A}_{1}$ ).

$$
\prod_{n>0}\left(1-q^{n}\right)\left(1-q^{n-1} z\right)\left(1-q^{n} z^{-1}\right)=\sum_{m \in \mathbb{Z}}(-1)^{m} z^{m} q^{m(m-1) / 2}
$$

This is a classical identity known as Jacobi's triple product identity.

As a second example we suppose $L$ has type $\tilde{A}_{1}^{\prime}$. Then $L^{0}$ has type $A_{1}$ and $\Phi^{0}=\left\{\alpha_{1},-\alpha_{1}\right\}$ as before, but we now have $h^{\mathrm{v}}=3$ and $M^{*}=\frac{1}{2} \mathbb{Z} \alpha_{1}$. We have $a_{0}=2, a_{1}=1, c_{0}=1, c_{1}=2$, thus

$$
\left\langle\alpha_{1}, \alpha_{1}\right\rangle=\frac{2 c_{1}}{a_{1}}=4 \quad \text { and } \quad \delta=2 \alpha_{0}+\alpha_{1} .
$$

We write $e_{-\alpha_{0}}=z$. Then $e_{-\alpha_{1}}=z^{-2} q$. The left-hand side of Macdonald's identity is

$$
\prod_{n>0}\left(1-q^{n}\right)\left(1-q^{n} z^{-1}\right)\left(1-q^{n-1} z\right)\left(1-q^{2 n+1} z^{-2}\right)\left(1-q^{2 n-1} z^{2}\right) .
$$

We have

$$
c\left(\frac{3}{2} n \alpha_{1}\right)=\left(\frac{9}{4} n^{2}+\frac{3}{2} n\right)\left\langle\alpha_{1}, \alpha_{1}\right\rangle=9 n^{2}+6 n .
$$

Also

$$
\chi^{0}\left(\frac{3}{2} n \alpha_{1}\right)=\frac{\left(z^{-2} q\right)^{-\frac{3}{2} n}-\left(z^{-2} q\right)^{\frac{3}{2} n+1}}{1-z^{-2} q}
$$

Thus the right-hand side of the identity is

$$
\begin{aligned}
\sum_{n \in \mathbb{Z}} & \frac{\left(\left(z^{-2} q\right)^{-\frac{3}{2} n}-\left(z^{-2} q\right)^{\frac{3}{2} n+1}\right)}{1-z^{-2} q} q^{\frac{n(3 n+2)}{2}} \\
& =\frac{1}{1-z^{-2} q}\left(\sum_{n \in \mathbb{Z}} z^{3 n} q^{\frac{n(3 n-1)}{2}}-\sum_{n \in \mathbb{Z}} z^{-3 n-2} q^{\frac{(n+1)(3 n+2)}{2}}\right) \\
& =\frac{1}{1-z^{-2} q}\left(\sum_{n \in \mathbb{Z}} z^{3 n} q^{\frac{n(3 n-1)}{2}}-\sum_{n \in \mathbb{Z}} z^{-3 n+1} q^{\frac{n(3 n-1)}{2}}\right) \\
& =\frac{1}{1-z^{-2} q} \sum_{n \in \mathbb{Z}}\left(z^{3 n}-z^{-3 n+1}\right) q^{\frac{n(3 n-1)}{2}} .
\end{aligned}
$$

We multiply both sides of the identity by $1-z^{-2} q$ and obtain
Proposition 20.6 (Macdonald's identity for type $\tilde{A}_{1}^{\prime}$ ).

$$
\begin{aligned}
\prod_{n>0} & \left(1-q^{n}\right)\left(1-q^{n} z^{-1}\right)\left(1-q^{n-1} z\right)\left(1-q^{2 n-1} z^{-2}\right)\left(1-q^{2 n-1} z^{2}\right) \\
& =\sum_{n \in \mathbb{Z}}\left(z^{3 n}-z^{-3 n+1}\right) q^{\frac{n(3 n-1)}{2}} .
\end{aligned}
$$

This is also a classical identity known as the quintuple product identity.

### 20.2 Specialisations of Macdonald's identities

We can obtain some striking identities, simpler than the original Macdonald identities, by specialising the latter identities in various ways. One way of specialising is simply to replace $e_{\alpha}$ by 1 for all $\alpha \in \Phi^{0}$. When this is done the expression $\chi^{0}(\lambda)$ is replaced by $d^{0}(\lambda)$ where

$$
d^{0}(\lambda)=\frac{\prod_{\alpha \in\left(\Phi^{0}\right)+}\left\langle\lambda+\rho^{0}, \alpha\right\rangle}{\prod_{\alpha \in\left(\Phi^{0}\right)^{+}}\left\langle\rho^{0}, \alpha\right\rangle} .
$$

This is shown in Theorem 12.19. The identities obtained by specialisation in this way involve Euler's $\phi$-function

$$
\phi(q)=(1-q)\left(1-q^{2}\right)\left(1-q^{3}\right) \ldots
$$

If we specialise the identity of Theorem 20.3 we obtain the following.

Theorem 20.7 (Macdonald's $\phi$-function identity).

$$
\phi(q)^{\operatorname{dim} L^{0}}=\sum_{\alpha \in M^{*}} d^{0}\left(h^{\mathrm{v}} \alpha\right) q^{c\left(h_{\alpha}^{\mathrm{v}}\right) / 2 h^{\mathrm{v}}}
$$

Proof. The left-hand side of the specialised identity is

$$
\phi(q)^{l+\left|\Phi^{0}\right|}=\phi(q)^{\operatorname{dim} L^{0}}
$$

On the right-hand side $\chi^{0}\left(h^{\mathrm{v}} \alpha\right)$ specialises to $d^{0}\left(h^{\mathrm{v}} \alpha\right)$.
We give some examples of this $\phi$-function identity.
$\underline{\text { Type } \tilde{A}_{1}}$

$$
\phi(q)^{3}=\sum_{n_{1} \in \mathbb{Z}}\left(4 n_{1}+1\right) q^{n_{1}\left(2 n_{1}+1\right)}
$$

Type $\tilde{A}_{2}$

$$
\begin{gathered}
\phi(q)^{8}=\sum_{\left(n_{1}, n_{2}\right) \in \mathbb{Z}^{2}} \frac{1}{2}\left(6 n_{1}-3 n_{2}+1\right)\left(-3 n_{1}+6 n_{2}+1\right)\left(3 n_{1}+3 n_{2}+2\right) \\
\times q^{3 n_{1}^{2}-3 n_{1} n_{2}+3 n_{2}^{2}+n_{1}+n_{2}} .
\end{gathered}
$$

$\underline{\text { Type } \tilde{C}_{2}}$

$$
\begin{aligned}
\phi(q)^{10}=\sum_{\left(n_{1}, n_{2}\right) \in \mathbb{Z}^{2}} & \left(12 n_{1}-6 n_{2}+1\right)\left(-6 n_{1}+6 n_{2}+1\right)\left(2 n_{2}+1\right)\left(3 n_{1}+1\right) \\
& \times q^{6 n_{1}^{2}-6 n_{1} n_{2}+3 n_{2}^{2}+n_{1}+n_{2}}
\end{aligned}
$$

## Type $\tilde{G}_{2}$

$$
\begin{aligned}
\phi(q)^{14}= & \sum_{\left(n_{1}, n_{2}\right) \in \mathbb{Z}^{2}} \frac{1}{15}\left(8 n_{1}-12 n_{2}+1\right)\left(-12 n_{1}+24 n_{2}+1\right)\left(3 n_{1}-3 n_{2}+1\right) \\
& \times\left(12 n_{2}+5\right)\left(-2 n_{1}+6 n_{2}+1\right)\left(4 n_{1}+3\right) \\
& \times q^{4 n_{1}^{2}-12 n_{1} n_{2}+12 n_{2}^{2}+n_{1}+n_{2}} .
\end{aligned}
$$

We next specialise the identities of Theorem 20.4 for twisted affine Kac-Moody algebras.

Theorem 20.8 (Macdonald's twisted $\phi$-function identities).
(a) The left-hand side of the identity is given as follows.
$\tilde{B}_{l}^{\mathrm{t}} \quad \phi(q)^{2 l^{2}-l-1} \phi\left(q^{2}\right)^{2 l+1}$
$\tilde{C}_{l}^{\mathrm{t}} \quad \phi(q)^{2 l+1} \phi\left(q^{2}\right)^{2 l^{2}-l-1}$
$\tilde{F}_{4}^{\mathrm{t}} \quad \phi(q)^{26} \phi\left(q^{2}\right)^{26}$
$\tilde{G}_{2}^{\mathrm{t}} \quad \phi(q)^{7} \phi\left(q^{3}\right)^{7}$
$\tilde{C}_{l}^{\prime} \quad \phi\left(q^{\frac{1}{2}}\right)^{2 l} \phi(q)^{2 l^{2}-3 l} \phi\left(q^{2}\right)^{2 l}$
$\tilde{A}_{1}^{\prime} \quad \phi\left(q^{\frac{1}{2}}\right)^{2} \phi(q)^{-1} \phi\left(q^{2}\right)^{2}$
(b) The right-hand side of the identity is

$$
\sum_{\alpha \in M^{*}} d^{0}\left(h^{\mathrm{v}} \alpha\right) q^{c\left(h^{\mathrm{v}} \alpha\right) / 2 h^{\mathrm{v}}}
$$

where $M^{*}$ is as in Theorem 20.4 (b).
We give some examples of twisted $\phi$-function identities.

## Type $\tilde{A}_{1}^{\prime}$

$$
\phi\left(q^{\frac{1}{2}}\right)^{2} \phi(q)^{-1} \phi\left(q^{2}\right)^{2}=\sum_{n_{1} \in \mathbb{Z}}\left(3 n_{1}+1\right) q^{\frac{1}{2} n_{1}\left(3 n_{1}+2\right)}
$$

## Type $\tilde{C}_{2}^{\mathrm{t}}$

$$
\begin{aligned}
\phi(q)^{5} \phi\left(q^{2}\right)^{5}=\sum_{\left(n_{1}, n_{2}\right) \in \mathbb{Z}^{2}} & \frac{1}{3}\left(8 n_{1}-4 n_{2}+1\right)\left(-8 n_{1}+8 n_{2}+1\right) \\
& \times\left(8 n_{1}+3\right)\left(2 n_{2}+1\right) \\
& \times q^{8 n_{1}^{2}-8 n_{1} n_{2}+4 n_{2}^{2}+2 n_{1}+n_{2}} .
\end{aligned}
$$

Type $\tilde{C}_{2}^{\prime}$

$$
\begin{aligned}
\phi\left(q^{\frac{1}{2}}\right)^{4} \phi(q)^{2} \phi\left(q^{2}\right)^{4}=\sum_{\left(n_{1}, n_{2}\right) \in \mathbb{Z}^{2}} & \frac{1}{6} \\
& \times\left(10 n_{1}-5 n_{2}+1\right)\left(-5 n_{1}+5 n_{2}+1\right) \\
& \times q^{5 n_{1}^{2}-5 n_{1} n_{2}+\frac{5}{2} n_{2}^{2}+n_{1}+n_{2}}
\end{aligned}
$$

## Type $\tilde{G}_{2}^{t}$

$$
\begin{aligned}
\phi(q)^{7} \phi\left(q^{3}\right)^{7}=\sum_{\left(n_{1}, n_{2}\right) \in \mathbb{Z}^{2}} & \frac{1}{10}\left(12 n_{1}-18 n_{2}+1\right)\left(-6 n_{1}+12 n_{2}+1\right) \\
& \times\left(-3 n_{1}+9 n_{2}+2\right)\left(6 n_{1}+5\right)\left(3 n_{1}-3 n_{2}+1\right) \\
& \times\left(2 n_{2}+1\right) q^{6 n_{1}^{2}-18 n_{1} n_{2}+18 n_{2}^{2}+n_{1}+3 n_{2}} .
\end{aligned}
$$

Another possibility to obtain specialised identities in one variable from Macdonald's identity is to apply a homomorphism

$$
\theta: \mathbb{C}\left[\left[e_{-\alpha_{0}}, e_{-\alpha_{1}}, \ldots, e_{-\alpha_{l}}\right]\right] \rightarrow \mathbb{C}[[q]]
$$

between rings of formal power series, given by

$$
\theta\left(e_{-\alpha_{0}}\right)=q^{s_{0}}, \theta\left(e_{-\alpha_{1}}\right)=q^{s_{1}}, \ldots, \theta\left(e_{-\alpha_{l}}\right)=q^{s_{l}}
$$

where $s_{0}, s_{1}, \ldots, s_{l}$ are non-negative integers. Of course under such a specialisation $e_{-\delta}$ would be mapped to $q^{a_{0} s_{0}+\cdots+a_{l} s_{l}}$, so that $q$ would have to be replaced by this power of $q$ in our earlier description of Macdonald's identity.

For example in type $\tilde{A}_{1}$ we obtain the following.
Proposition 20.9 (Macdonald's 1-variable identity for $\tilde{A}_{1}$ ).

$$
\begin{aligned}
\prod_{n>0} & \left(1-q^{\left(s_{0}+s_{1}\right) n}\right)\left(1-q^{s_{0}(n-1)+s_{1} n}\right)\left(1-q^{s_{0} n+s_{1}(n-1)}\right) \\
& =\sum_{m \in \mathbb{Z}}(-1)^{m} q^{s_{0} \frac{m(m-1)}{2}+s_{1} \frac{m(m+1)}{2}} .
\end{aligned}
$$

We mention some explicit examples of this identity. If $\left(s_{0}, s_{1}\right)=(1,1)$ we obtain

$$
\frac{\phi(q)^{2}}{\phi\left(q^{2}\right)}=\sum_{m \in \mathbb{Z}}(-1)^{m} q^{m^{2}}
$$

that is

$$
\begin{aligned}
& (1-q)^{2}\left(1-q^{2}\right)\left(1-q^{3}\right)^{2}\left(1-q^{4}\right)\left(1-q^{5}\right)^{2}\left(1-q^{6}\right) \cdots \\
& \quad=1-2 q+2 q^{4}-2 q^{9}+2 q^{16}-\cdots
\end{aligned}
$$

This is a classical formula of Gauss.
Next consider the example given by $\left(s_{0}, s_{1}\right)=(2,1)$. Then we obtain

$$
\phi(q)=\sum_{m \in \mathbb{Z}}(-1)^{m} q^{\frac{m(3 m-1)}{2}}
$$

that is

$$
\begin{aligned}
& (1-q)\left(1-q^{2}\right)\left(1-q^{3}\right)\left(1-q^{4}\right)\left(1-q^{5}\right)\left(1-q^{6}\right) \cdots \\
& \quad=1-q-q^{2}+q^{5}+q^{7}-q^{12}-q^{15}+q^{22}+q^{26}-\cdots .
\end{aligned}
$$

This is a well known formula of Euler.
Many additional formulae can be obtained by taking different values of $\left(s_{0}, s_{1}\right)$ or different affine Kac-Moody algebras.

### 20.3 Irreducible modules for affine algebras

We next consider the weights of the irreducible modules $L(\lambda), \lambda \in X^{+}$, for the affine Kac-Moody algebra $L(A)$. We recall that $X$ is the set of $\lambda \in H^{*}$ with $\lambda\left(h_{i}\right) \in \mathbb{Z}$ for $i=0,1, \ldots, l$ and $X^{+}$is the set of $\lambda \in X$ with $\lambda\left(h_{i}\right) \geq 0$ for $i=0,1, \ldots, l$.

It is convenient to introduce the fundamental weights $\omega_{0}, \omega_{1}, \ldots, \omega_{l} . \omega_{i}$ is the element of $X^{+}$defined by

$$
\omega_{i}\left(h_{j}\right)=\delta_{i j} \quad \omega_{i}(d)=0 .
$$

Since the imaginary root $\delta$ satisfies

$$
\delta\left(h_{j}\right)=0 \quad \delta(d)=1
$$

we see that $\omega_{0}, \omega_{1}, \ldots, \omega_{l}, \delta$ form a basis of $H^{*}$.
If $\lambda \in H^{*}$ satisfies

$$
\lambda=\xi_{0} \omega_{0}+\xi_{1} \omega_{1}+\cdots+\xi_{l} \omega_{l}+\xi \delta
$$

then $\lambda$ lies in $X$ if and only if $\xi_{i} \in \mathbb{Z}$ for $i=0,1, \ldots, l . \xi$ can be any element of $\mathbb{C}$. Also $\lambda \in X^{+}$if and only if $\xi_{i} \in \mathbb{Z}$ and $\xi_{i} \geq 0$ for $i=0,1, \ldots, l$.

Now every weight $\mu$ of $L(\lambda)$ has form $\mu=\lambda-m_{0} \alpha_{0}-m_{1} \alpha_{1}-\cdots-m_{l} \alpha_{l}$ for certain $m_{i} \in \mathbb{Z}$ with $m_{i} \geq 0$. Since $\alpha_{i}(c)=0$ for $i=0,1, \ldots, l$ we have $\mu(c)=\lambda(c)$. Thus all the weights $\mu$ of $L(\lambda)$ have the same value of $\mu(c)$. Since $c=c_{0} h_{0}+c_{1} h_{1}+\cdots+c_{l} h_{l}$ we have $\lambda(c)=\sum_{i=0} c_{i} \lambda\left(h_{i}\right)$ and so, for $\lambda \in X^{+}, \lambda(c)$ is a non-negative integer. The integer $\lambda(c)$ is called the level of the module $L(\lambda)$.

Proposition 20.10 If $L(\lambda)$ has level 0 then $\lambda=\xi \delta$ for some $\xi \in \mathbb{C}$ and $\operatorname{dim} L(\lambda)=1$.

Proof. If $\lambda(c)=0$ then $\lambda\left(h_{i}\right)=0$ for $i=0,1, \ldots, l$. Writing $\lambda=\xi_{0} \omega_{0}+\cdots+$ $\xi_{l} \omega_{l}+\xi \delta$ we see that $\xi_{i}=0$ for $i=0,1, \ldots, l$, hence $\lambda=\xi \delta$. Kac' character formula then shows that $\operatorname{ch} L(\xi \delta)=e_{\xi \delta}$. Thus $\operatorname{dim} L(\xi \delta)=1$.

Since the modules $L(\lambda)$ of level 0 are trivial 1-dimensional modules we shall subsequently concentrate on modules $L(\lambda), \lambda \in X^{+}$, of level greater than 0 .

If $\mu$ is a weight of $L(\lambda)$ then so is $w(\mu)$ for any $w \in W$, by Proposition 19.13. We take a $w \in W$ for which the height of $\lambda-w(\mu)$ is minimal and put $\nu=w(\mu)$. Since $s_{i}(\nu)=\nu-\nu\left(h_{i}\right) \alpha_{i}$ the minimality of the height shows that $\nu\left(h_{i}\right) \geq 0$. Thus for any weight $\mu$ of $L(\lambda)$ there exists $w \in W$ with $w(\mu) \in X^{+}$. We now show that the converse is true also.

Theorem 20.11 Let $\lambda \in X^{+}$have $\lambda(c)>0$. Then $\mu \in X$ is a weight of $L(\lambda)$ if and only if there exists $w \in W$ such that $w(\mu) \in X^{+}$and $w(\mu) \prec \lambda$.

Proof. It will be sufficient to show that if $\mu \in X^{+}$with $\mu \prec \lambda$ then $\mu$ is a weight of $L(\lambda)$. The proof of this is non-trivial and reminiscent of that of Proposition 16.23.

Let $\mu=\lambda-\alpha$ where $\alpha=\sum k_{i} \alpha_{i}$ and each $k_{i} \geq 0$. We may assume $k_{i}>0$ for some $i$. supp $\alpha$ is the set of $i$ for which $k_{i}>0$. We first show that every connected component of $\operatorname{supp} \alpha$ contains an $i$ with $\lambda\left(h_{i}\right)>0$. Suppose if possible there exists a connected component $S$ of $\operatorname{supp} \alpha$ with $\lambda\left(h_{i}\right)=0$ for all $i \in S$. We have

$$
L(\lambda)_{\mu} \subset \mathfrak{H}\left(N^{-}\right)_{-\alpha} v_{\lambda}
$$

where $v_{\lambda}$ is a highest weight vector of $L(\lambda)$ and, by the PBW basis theorem, $\mathfrak{U}\left(N^{-}\right)_{-\alpha}$ is spanned by elements of the form

$$
\prod_{\beta \in \Phi^{+}} e_{-\beta}^{k_{\beta}}
$$

where $k_{\beta} \geq 0, \sum k_{\beta} \beta=\alpha$, and each $\beta$ involves fundamental roots which all lie in the same connected component of $\operatorname{supp} \alpha$. (We recall from Proposition 16.21 that $\operatorname{supp} \beta$ is connected.) Now the $e_{-\beta}$ with fundamental roots in different connected components of $\operatorname{supp} \alpha$ commute with one another, so we may bring the $e_{-\beta}$ with fundamental roots in $S$ to the right of the above product. But for such $\beta$ we have $e_{-\beta} v_{\lambda}=0$.

For $f_{i} v_{\lambda}=0$ for each $i \in S$ by Theorem 19.19, since $\lambda\left(h_{i}\right)=0$. It follows that

$$
\mathfrak{U}\left(N^{-}\right)_{-\alpha} v_{\lambda}=O
$$

and so $L(\lambda)_{\mu}=O$, a contradiction. Hence there exists $i \in S$ with $\lambda\left(h_{i}\right)>0$.
Now let $\Psi$ be defined by

$$
\Psi=\left\{\gamma \in Q^{+} ; \gamma \prec \alpha, \lambda-\gamma \text { is a weight of } L(\lambda)\right\}
$$

The set $\Psi$ is finite. Let $\beta \in \Psi$ be an element of maximal height. Then $\beta \prec \alpha$. We aim to show that $\beta=\alpha$ and hence that $\lambda-\alpha$ is a weight of $L(\lambda)$. Let $\beta=\sum m_{i} \alpha_{i}$ with each $m_{i} \geq 0$. We have $\alpha=\sum k_{i} \alpha_{i}$ with $m_{i} \leq k_{i}$ for each $i$.

Let $I=\{0,1, \ldots, l\}$ and $J$ be the subset of $I$ given by $J=\left\{i \in I ; k_{i}=m_{i}\right\}$. We aim to show that $J=I$ and so that $\beta=\alpha$. Suppose if possible that $J \neq I$. Consider the non-empty subset of $I$ given by $\operatorname{supp} \alpha-(\operatorname{supp} \alpha \cap J)$. This set splits into connected components. Let $M$ be a connected component of $\operatorname{supp} \alpha-(\operatorname{supp} \alpha \cap J)$. Let $i \in M$. Then $\lambda-\beta$ is a weight of $L(\lambda)$ but $\lambda-\beta-\alpha_{i}$ is not. Thus $(\lambda-\beta)\left(h_{i}\right) \leq 0$. Also $\mu\left(h_{i}\right) \geq 0$ since $\mu \in X^{+}$and so $(\lambda-\alpha)\left(h_{i}\right) \geq 0$. Thus we have

$$
\alpha\left(h_{i}\right) \leq \lambda\left(h_{i}\right) \leq \beta\left(h_{i}\right) .
$$

Let $\gamma=\sum_{j \in M}\left(k_{j}-m_{j}\right) \alpha_{j}$. We have $k_{j}-m_{j}>0$ for all $j \in M$. We also have

$$
\gamma\left(h_{i}\right)=\sum_{j \in M}\left(k_{j}-m_{j}\right) A_{i j}
$$

However, $\gamma\left(h_{i}\right)=(\alpha-\beta)\left(h_{i}\right)$ since $\operatorname{supp}(\alpha-\beta)=\operatorname{supp} \alpha-J$ and $M$ is a connected component of supp $\alpha-J$. Thus $\gamma\left(h_{i}\right) \leq 0$ for each $i \in M$.

Let $A_{M}$ be the principal minor $\left(A_{i j}\right)$ for $i, j \in M$. Let $u$ be the column vector with entries $k_{i}-m_{i}$ for $i \in M$. Then we have $u>0$ and $A_{M} u \leq 0$. If $M$ has finite type $A_{M}(-u) \geq 0$ would imply $-u>0$ or $-u=0$. Thus $M$ does not have finite type. Since $M$ is a subset of $I$ which has affine type we must have $M=I$ by Lemma 15.13. Thus $\operatorname{supp} \alpha=I$ and $J=\phi$. But then, for all $i \in I, \lambda-\beta$ is a weight of $L(\lambda)$ but $\lambda-\beta-\alpha_{i}$ is not. Thus $(\lambda-\beta)\left(h_{i}\right) \leq 0$ for all $i \in I$. Hence $\alpha\left(h_{i}\right) \leq \lambda\left(h_{i}\right) \leq \beta\left(h_{i}\right)$ for all $i \in I$. We now have $u>0$ and $A u \leq 0$. Since $A$ is affine we can deduce $A u=0$. This shows that $\alpha\left(h_{i}\right)=\beta\left(h_{i}\right)$ for all $i \in I$. Hence $\alpha\left(h_{i}\right)=\lambda\left(h_{i}\right)$ for all $i \in I$, that is $\mu\left(h_{i}\right)=0$ for each $i$. But then we have $\mu(c)=0$, and so $\lambda(c)=0$, a contradiction.

Corollary 20.12 If $\mu$ is a weight of $L(\lambda)$ then $\mu-\delta$ is also a weight.

Proof. Since $\mu$ is a weight there exists $w \in W$ such that $w(\mu) \in X^{+}$. Then $w(\mu-\delta)=w(\mu)-\delta \in X^{+}$. Since $w(\mu)-\delta \prec \lambda$ it follows from Theorem 20.11 that $w(\mu)-\delta$ is a weight of $L(\lambda)$. Hence $\mu-\delta$ is also a weight.

It follows from this corollary that $\mu-i \delta$ is a weight for all positive integers $i$. On the other hand there exist only finitely many positive integers $i$ such that $\mu+i \delta \prec \lambda$.

Definition $A$ weight $\mu$ of $L(\lambda)$ is called a maximal weight if $\mu+\delta$ is not a weight.

Corollary 20.13 For each weight $\mu$ of $L(\lambda)$ there are a unique maximal weight $\nu$ and a unique non-negative integer $i$ such that $\mu=\nu-i \delta$.

Proof. Consider the sequence $\mu, \mu+\delta, \mu+2 \delta, \ldots$ There exists $i$ such that $\mu+i \delta$ is a weight of $L(\lambda)$ but $\mu+(i+1) \delta$ is not a weight. Let $\nu=\mu+i \delta$. Then $\nu$ is a maximal weight of $L(\lambda)$ and $\mu=\nu-i \delta$.

If $\mu=\nu^{\prime}-i^{\prime} \delta$ where $\nu^{\prime}$ is a maximal weight and $i^{\prime}$ a non-negative integer we show $\nu=\nu^{\prime}$ and $i=i^{\prime}$. Otherwise we may assume $i<i^{\prime}$. Then $\nu^{\prime}=\nu+\left(i^{\prime}-i\right) \delta$ is a weight. By Corollary $20.12 \nu+\delta$ is also a weight. Thus $\nu$ is not a maximal weight and we have a contradiction.

A string of weights of $L(\lambda)$ is a set

$$
\nu, \nu-\delta, \nu-2 \delta, \ldots
$$

where $\nu$ is a maximal weight. Each weight lies in a unique string of weights. Thus it is natural to consider the set of maximal weights of $L(\lambda)$.

Proposition 20.14 The set of maximal weights of $L(\lambda), \lambda \in X^{+}$, is invariant under the Weyl group.

Proof. Let $w \in W$. Then $\mu$ is a weight if and only if $w(\mu)$ is a weight. Thus if $\mu$ is a maximal weight $w(\mu)$ is a weight but $w(\mu)+\delta=w(\mu+\delta)$ is not a weight. Thus $w(\mu)$ is a maximal weight.

Corollary 20.15 Each maximal weight of $L(\lambda), \lambda \in X^{+}$, has form $w(\mu)$ where $w \in W$ and $\mu$ is a dominant maximal weight.

We shall therefore consider the set of dominant maximal weights of $L(\lambda)$. We shall show that $L(\lambda)$ has only finitely many dominant maximal weights.

We recall from Section 17.3 that the fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$ was defined by

$$
\begin{array}{rlrl}
A^{*} & =\left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0\right. & \text { for } i=1, \ldots, l ; & \left.\lambda\left(h_{\theta}\right)<\frac{1}{a_{0}}\right\} \\
& =\left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ;\left\langle\lambda, \alpha_{i}\right\rangle>0\right. & \text { for } i=1, \ldots, l ; \quad\langle\lambda, \theta\rangle<1\} .
\end{array}
$$

Its closure $\overline{A^{*}}$ is a fundamental region for the action of $W$ on $\left(H_{\mathbb{R}}^{0}\right)^{*}$.
We also recall that

$$
H^{*}=\left(H^{0}\right)^{*} \oplus(\mathbb{C} \gamma+\mathbb{C} \delta)
$$

where $\left(H^{0}\right)^{*}$ is embedded in $H^{*}$ by assuming $\lambda(c)=0, \lambda(d)=0$ for $\lambda \in\left(H^{0}\right)^{*}$. By Lemma 20.1 we have, for $\lambda \in H^{*}$,

$$
\lambda=\lambda^{0}+\lambda(c) \gamma+a_{0}^{-1} \lambda(d) \delta
$$

where $\lambda^{0} \in\left(H^{0}\right)^{*}$. Let $Q^{0} \subset\left(H^{0}\right)^{*}$ be the set of $\lambda^{0}$ given by $\lambda$ in the root lattice $Q \subset H^{*}$.

Proposition 20.16 Let $\lambda \in X^{+}$have level $\lambda(c)=k>0$. Then the map $\mu \rightarrow \mu^{0}$ gives a bijection between the set of dominant maximal weights of $L(\lambda)$ and $\left(\lambda^{0}+Q^{0}\right) \cap \overline{k A^{*}}$.

Proof. Let $\mu$ be a dominant maximal weight of $L(\lambda)$. Then $\mu=\lambda-$ $m_{0} \alpha_{0}-\cdots-m_{l} \alpha_{l}$ for certain $m_{i} \in \mathbb{Z}$ with $m_{i} \geq 0$. Hence $\mu^{0}=\lambda^{0}-$ $\left(m_{0} \alpha_{0}+\cdots+m_{l} \alpha_{l}\right)^{0}$ and so $\mu^{0} \in \lambda^{0}+Q^{0}$.

Now $\mu=\mu^{0}+k \gamma+a_{0}^{-1} \mu(d) \delta$. Since $\mu \in X^{+}$we have $\mu\left(h_{i}\right) \geq 0$ for $i=0$, $1, \ldots, l$. Now $\gamma\left(h_{i}\right)=\delta\left(h_{i}\right)=0$ for $i=1, \ldots, l$ and so $\mu^{0}\left(h_{i}\right) \geq 0$ for $i=$ $1, \ldots, l$. We also have

$$
\left\langle\mu^{0}, \theta\right\rangle=\langle\mu, \theta\rangle=\left\langle\mu, \delta-a_{0} \alpha_{0}\right\rangle=\mu(c)-\mu\left(h_{0}\right)=k-\mu\left(h_{0}\right)
$$

Since $\mu\left(h_{0}\right) \geq 0$ we have $\left\langle\mu^{0}, \theta\right\rangle \leq k$. Thus $\mu^{0} \in \overline{k A^{*}}$.
Hence $\mu \rightarrow \mu^{0}$ maps dominant maximal weights of $L(\lambda)$ into $\left(\lambda^{0}+Q^{0}\right) \cap$ $\overline{k A^{*}}$. We wish to show this map is bijective. We first show it is surjective. Let $\nu \in\left(\lambda^{0}+Q^{0}\right) \cap \overline{k A^{*}}$. Then, since $\alpha_{i}^{0}=\alpha_{i}$ for $i=1, \ldots, l$ and

$$
\alpha_{0}^{0}=\left(-a_{0}^{-1} \theta+a_{0}^{-1} \delta\right)^{0}=-a_{0}^{-1} \theta
$$

we have

$$
\nu=\lambda^{0}+k_{1} \alpha_{1}+\cdots+k_{l} \alpha_{l}-k_{0} a_{0}^{-1} \theta
$$

for certain $k_{0}, k_{1}, \ldots, k_{l} \in \mathbb{Z}$. Since $\theta=a_{1} \alpha_{1}+\cdots+a_{l} \alpha_{l}$ we have

$$
\nu=\lambda^{0}+\left(m-k_{0} a_{0}^{-1}\right) \theta-\left(m a_{1}-k_{1}\right) \alpha_{1}-\cdots-\left(m a_{l}-k_{l}\right) \alpha_{l} .
$$

We choose $m \in \mathbb{Z}$ with $m \geq k_{i} / a_{i}$ for $i=0,1, \ldots, l$. Then

$$
\nu=\lambda^{0}+m_{0}\left(a_{0}^{-1} \theta\right)-m_{1} \alpha_{1}-\cdots-m_{l} \alpha_{l}
$$

where $m_{i}=m a_{i}-k_{i}$ for $i=0,1, \ldots, l$. Thus the $m_{i}$ are non-negative integers for $i=0,1, \ldots, l$. Let $\mu=\lambda-m_{0} \alpha_{0}-\cdots-m_{l} \alpha_{l}$. Then

$$
\mu^{0}=\lambda^{0}+m_{0}\left(a_{0}^{-1} \theta\right)-m_{1} \alpha_{1}-\cdots-m_{l} \alpha_{l}=\nu
$$

We show that $\mu \in X^{+}$. We have $\mu\left(h_{i}\right)=\mu^{0}\left(h_{i}\right)=\nu\left(h_{i}\right) \geq 0$ for $i=1, \ldots, l$. Also $\mu\left(h_{0}\right)=k-\left\langle\mu^{0}, \theta\right\rangle=k-\langle\nu, \theta\rangle \geq 0$. Hence $\mu \in X^{+}$and $\mu \prec \lambda$. Thus $\mu$ is a dominant weight of $L(\lambda)$ by Theorem 20.11. Hence we have shown that $\nu=\mu^{0}$ for some dominant weight $\mu$ of $L(\lambda)$. By replacing $\mu$ by the maximal weight in the chain of weights containing $\mu$ we may assume that $\mu$ is a dominant maximal weight. Thus our map is surjective.

To show the map is injective let $\mu, \mu^{\prime}$ be dominant maximal weights of $L(\lambda)$ with $\mu^{0}=\left(\mu^{\prime}\right)^{0}$. We have

$$
\begin{aligned}
& \mu=\mu^{0}+k \gamma+a_{0}^{-1} \mu(d) \delta \\
& \mu^{\prime}=\left(\mu^{\prime}\right)^{0}+k \gamma+a_{0}^{-1} \mu^{\prime}(d) \delta
\end{aligned}
$$

hence $\mu-\mu^{\prime}=a_{0}^{-1}\left(\mu(d)-\mu^{\prime}(d)\right) \delta$. Now $\lambda-\mu \in Q$ and $\lambda-\mu^{\prime} \in Q$ hence $\mu-\mu^{\prime} \in Q$ and $a_{0}^{-1}\left(\mu(d)-\mu^{\prime}(d)\right) \delta \in Q$. This shows that $a_{0}^{-1}\left(\mu(d)-\mu^{\prime}(d)\right) \in$ $\mathbb{Z}$. Thus $\mu=\mu^{\prime}+r \delta$ for some $r \in \mathbb{Z}$. Since $\mu, \mu^{\prime}$ are both maximal weights we must have $r=0$. Thus $\mu^{\prime}=\mu$.

Corollary 20.17 The set of dominant maximal weights of $L(\lambda), \lambda \in X^{+}$, is finite.

Proof. $Q^{0}$ is a lattice in $\left(H^{0}\right)^{*}$, that is a free abelian subgroup whose rank is the dimension of $\left(H^{0}\right)^{*} . \lambda^{0}+Q^{0}$ is a coset of this lattice. On the other hand the set $\overline{k A^{*}}$ is bounded. Hence the intersection $\left(\lambda^{0}+Q^{0}\right) \cap \overline{k A^{*}}$ must be finite. Thus the set of dominant maximal weights is also finite, by Proposition 20.16.

We now have a procedure for describing all weights of $L(\lambda), \lambda \in X^{+}$. First determine the finite set $\left(\lambda^{0}+Q^{0}\right) \cap \overline{k A^{*}}$ where $k=\lambda(c)$. For each element $\nu$ in this finite set there is a unique dominant maximal weight $\mu$ of $L(\lambda)$ with $\mu^{0}=\nu$. This gives the set of all dominant maximal weights. By applying elements of the Weyl group to these we obtain all maximal weights. Finally
by subtracting positive integral multiples of $\delta$ from the maximal weights we obtain all weights of $L(\lambda)$.

We next consider the weights in a string

$$
\mu, \mu-\delta, \mu-2 \delta, \ldots
$$

We wish to show that the multiplicities of these weights form an increasing function as we move down the string, i.e. that $m_{\mu-(i+1) \delta} \geq m_{\mu-i \delta}$ for all $i \geq 0$. In order to do this we consider $L(\lambda)$ as a $T$-module where $T$ is the subalgebra of $L(A)$ given by

$$
T=\cdots \oplus L_{-2 \delta} \oplus L_{-\delta} \oplus H \oplus L_{\delta} \oplus L_{2 \delta} \oplus \cdots
$$

Thus $T$ is spanned by $H$ and the root spaces for the imaginary roots. The algebra $T$ has a triangular decomposition

$$
T=T^{-} \oplus H \oplus T^{+}
$$

where $T^{-}=\sum_{i>0} L_{-i \delta}, T^{+}=\sum_{i>0} L_{i \delta}$. One can define the category $\mathcal{O}$ of Tmodules in a manner analogous to that in Section 19.1. One can also define Verma modules for $T$. If $\lambda \in H^{*}$ we define

$$
M(\lambda)=\mathfrak{U}(T) / \mathfrak{H}(T) T^{+}+\sum_{x \in H} \mathfrak{H}(T)(x-\lambda(x))
$$

This is the Verma module for $T$ with highest weight $\lambda$. There is a bijection $\mathfrak{U}\left(T^{-}\right) \rightarrow M(\lambda)$ given by $u \rightarrow u m_{\lambda}$ where $m_{\lambda} \in M(\lambda)$ is the image of $1 \in \mathfrak{H}(T)$.

We shall investigate properties of Verma modules for $T$ by considering the expression

$$
\Omega_{0}=2 \sum_{i>0} \sum_{j} e_{-i \delta}^{(j)} e_{i \delta}^{(j)}
$$

where $e_{i \delta}^{(j)}$ is a basis for $L_{i \delta}$ and $e_{-i \delta}^{(j)}$ is the dual basis for $L_{-i \delta}$. Thus

$$
\left\langle e_{i \delta}^{(j)}, e_{-i \delta}^{(k)}\right\rangle=\delta_{j k}
$$

and $\left[e_{i \delta}^{(j)}, e_{-i \delta}^{(k)}\right]=\delta_{j k} i c$ by Corollary 16.5.
Although the expression for $\Omega_{0}$ is an infinite sum the action of $\Omega_{0}$ on any $T$-module in category $\mathcal{O}$ is well defined, since all but a finite number of the terms will act as zero.

Lemma 20.18 Let $\lambda \in H^{*}$ and $M(\lambda)$ be the associated Verma module for $T$. Let $u \in \mathfrak{U}(T)_{m \delta}$ where $m \in \mathbb{Z}$ and $m \neq 0$. Then $\Omega_{0} u-u \Omega_{0}$ acts on $M(\lambda)$ in the same way as $-2 \lambda(c) m u$.

Proof. $u$ is a linear combination of products of elements, each in $T_{r \delta}$ for some $r$ with $r \neq 0$.

First suppose $u \in T_{r \delta}$. We assume that $u$ is one of the basis elements $u=e_{r \delta}^{(j)}$. Then $u$ commutes with all $e_{i \delta}^{(k)}, e_{-i \delta}^{(k)}$ except for $e_{-r \delta}^{(j)}$. Thus

$$
\begin{aligned}
\Omega_{0} u-u \Omega_{0} & =2\left(e_{-r \delta}^{(j)} e_{r \delta}^{(j)} e_{r \delta}^{(j)}-e_{r \delta}^{(j)} e_{-r \delta}^{(j)} e_{r \delta}^{(j)}\right) \\
& =-2 r c e_{r \delta}^{(j)}=-2 r u c=-2 r \lambda(c) u
\end{aligned}
$$

on $M(\lambda)$. The same will then apply to any $u \in T_{r \delta}$.
Next suppose $u=u_{1} u_{2}$ where

$$
\begin{aligned}
& \Omega_{0} u_{1}-u_{1} \Omega_{0}=-2 \lambda(c) r_{1} u_{1} \\
& \Omega_{0} u_{2}-u_{2} \Omega_{0}=-2 \lambda(c) r_{2} u_{2} \quad \text { on } M(\lambda)
\end{aligned}
$$

Then

$$
\begin{aligned}
\Omega_{0} u-u \Omega_{0} & =\Omega_{0} u_{1} u_{2}-u_{1} u_{2} \Omega_{0} \\
& =u_{1} \Omega_{0} u_{2}-2 \lambda(c) r_{1} u-u_{1} \Omega_{0} u_{2}-2 \lambda(c) r_{2} u \\
& =-2 \lambda(c)\left(r_{1}+r_{2}\right) u \quad \text { on } M(\lambda)
\end{aligned}
$$

The required result then follows for arbitrary $u \in \mathfrak{U}(T)_{m \delta}$ by taking linear combinations of such repeated products.

Proposition 20.19 Let $\lambda \in H^{*}$ satisfy $\lambda(c)>0$. Then the Verma module $M(\lambda)$ for $T$ is irreducible.

Proof. Suppose if possible that $M(\lambda)$ has a proper submodule $K$. Let $v$ be a highest weight vector of $K$. Then $v \in M(\lambda)_{\lambda-m \delta}$ for some $m \in \mathbb{Z}$ with $m>0$. Thus $v=u m_{\lambda}$ for some $u \in \mathfrak{H}\left(T^{-}\right)_{-m \delta}$. We consider the actions

$$
\Omega_{0}: M(\lambda) \rightarrow M(\lambda) \quad u: M(\lambda) \rightarrow M(\lambda)
$$

By Lemma 20.18 we have

$$
\left(\Omega_{0} u-u \Omega_{0}\right) m_{\lambda}=2 \lambda(c) m u m_{\lambda} .
$$

Thus $\Omega_{0} v-u\left(\Omega_{0} m_{\lambda}\right)=2 \lambda(c) m v$. Now $\Omega_{0} m_{\lambda}=0$ and $\Omega_{0} v=0$ since $m_{\lambda}$ and $v$ are highest weight vectors in $M(\lambda)$ and $K$ respectively. Thus $2 \lambda(c) m v=0$. But $v \neq 0, m>0, \lambda(c)>0$ and so we have a contradiction. Thus $M(\lambda)$ is irreducible.

We now consider the structure of $L(\lambda)$ as a $T$-module.

Proposition 20.20 Suppose $\lambda \in X^{+}$with $\lambda(c)>0$. Then the $T$-module $L(\lambda)$ is completely reducible. Its irreducible components are Verma modules for $T$.

Proof. Let $U$ be the subspace of $L(\lambda)$ given by

$$
U=\left\{v \in L(\lambda) ; T^{+} v=0\right\}
$$

Let $B$ be a basis of $U$. We may choose $B$ to be a basis of weight vectors of $U$, i.e. so that each element of $B$ lies in a weight space $L(\lambda)_{\mu}$. Suppose $v \in B$ has weight $\mu$. Then $T^{+} v=0$ and $x v=\mu(x) v$ for $x \in H$, hence $T v=$ $T^{-} v$. Let $M(\mu)$ be the Verma module for $T$ with highest weight $\mu$. Then we have a homomorphism of $T$-modules $M(\mu) \rightarrow T v$ given by $u m_{\mu} \rightarrow u v$ for $u \in T^{-}$. Now $\mu=\lambda-i \delta$ for some $i \geq 0$ hence $\mu(c)=\lambda(c)>0$. Thus the Verma module $M(\mu)$ for $T$ is irreducible by Proposition 20.19. Hence the homomorphism $M(\mu) \rightarrow T v$ is an isomorphism and so $T v$ is a Verma module for $T$. Let $V=\sum_{v \in B} T v$. We claim that this sum of $T$-modules is a direct sum. For consider

$$
T v \cap \sum_{\substack{v^{\prime} \in B \\ v^{\prime} \neq v}} T v^{\prime}
$$

Since the Verma module $T v$ is irreducible we have $T v \cap U=\mathbb{C} v$. We also have

$$
\left(\sum_{\substack{v^{\prime} \in B \\ v^{\prime} \neq v}} T v^{\prime}\right) \cap U=\sum_{v^{\prime} \neq v} \mathbb{C} v^{\prime}
$$

Since $v \notin \sum_{v^{\prime} \neq v} \mathbb{C} v^{\prime}$ we see that $T v$ is not contained in $\sum_{v^{\prime} \neq v} T v^{\prime}$. Again, since $T v$ is irreducible we have $T v \cap \sum_{v^{\prime} \neq v} T v^{\prime}=O$. Hence $V=\bigoplus_{v \in B} T v$. Thus $V$ is a direct sum of Verma modules for $T$.

We wish to show that $V=L(\lambda)$. We suppose if possible that $V \neq L(\lambda)$. We consider the $T$-module $L(\lambda) / V$. Since $L(\lambda)=\bigoplus_{\mu} L(\lambda)_{\mu}$ and $V=\bigoplus_{\mu} V_{\mu}$ we have $L(\lambda) / V=\bigoplus_{\mu}(L(\lambda) / V)_{\mu}$. As $L(\lambda) / V$ is assumed to be non-zero we can find a weight $\mu$ of $L(\lambda) / V$ such that $\mu+i \delta$ is not a weight for any $i>0$. Then $T^{+}(L(\lambda) / V)_{\mu}=O$, that is $T^{+} L(\lambda)_{\mu} \subset V$.

We now consider the map $\Omega_{0}: L(\lambda) \rightarrow L(\lambda)$. Since the action of $\Omega_{0}$ preserves weight spaces we have $\Omega_{0}: L(\lambda)_{\mu} \rightarrow L(\lambda)_{\mu}$. The weight space
$L(\lambda)_{\mu}$ is finite dimensional, so decomposes into a direct sum of generalised eigenspaces of $\Omega_{0}$, given by

$$
L(\lambda)_{\mu}=\bigoplus_{\zeta \in \mathbb{C}}\left(L(\lambda)_{\mu}\right)_{\zeta}
$$

where $\left(\Omega_{0}-\zeta 1\right)^{k}=0$ on $\left(L(\lambda)_{\mu}\right)_{\zeta}$ for some $k$. Since $L(\lambda)_{\mu}$ does not lie in $V$ there exists $\zeta \in \mathbb{C}$ such that $\left(L(\lambda)_{\mu}\right)_{\zeta}$ does not lie in $V$. We choose $v \in\left(L(\lambda)_{\mu}\right)_{\zeta}$ with $v \notin V$. Then

$$
\left(\Omega_{0}-\zeta 1\right)^{k} v=0
$$

and $\Omega_{0} v \in V$, since $T^{+} L(\lambda)_{\mu} \subset V$. If $\zeta \neq 0$ the polynomials $(t-\zeta)^{k}$ and $t$ are coprime so we could deduce $v \in V$, a contradiction. Hence $\zeta=0$ and $\Omega_{0}^{k} v=0$.

Now $T^{+} v \neq 0$ since $v \notin V$. So there exist $m>0$ and $u \in \mathfrak{H}\left(T^{+}\right)_{m \delta}$ with $u v \neq 0$ and $T^{+}(u v)=0$. Let $v^{\prime}=u v$. Then $v^{\prime} \neq 0$ and $\Omega_{0} v^{\prime}=0$.

Now all the weights $\nu$ of $L(\lambda)$ satisfy $\nu(c)=\lambda(c)$. Thus we may apply the argument of Lemma 20.18 to $L(\lambda)$ and obtain

$$
\Omega_{0} u-u \Omega_{0}=-2 \lambda(c) m u \quad \text { on } L(\lambda)
$$

Then

$$
\Omega_{0} u v-u \Omega_{0} v=-2 \lambda(c) m u v
$$

that is

$$
\left(\Omega_{0}+2 \lambda(c) m\right) v^{\prime}=u \Omega_{0} v
$$

It follows that

$$
\left(\Omega_{0}+2 \lambda(c) m\right)^{2} v^{\prime}=\left(\Omega_{0}+2 \lambda(c) m\right)\left(u \Omega_{0} v\right)=u\left(\Omega_{0}^{2} v\right)
$$

and continuing thus we obtain

$$
\left(\Omega_{0}+2 \lambda(c) m\right)^{k} v^{\prime}=u\left(\Omega_{0}^{k} v\right)=0
$$

But $\lambda(c)>0$ and $m>0$, thus the polynomials $(t+2 \lambda(c) m)^{k}$ and $t$ are coprime. Thus $\left(\Omega_{0}+2 \lambda(c) m\right)^{k} v^{\prime}=0$ and $\Omega_{0} v^{\prime}=0$ imply $v^{\prime}=0$, a contradiction.

Thus we have obtained our required contradiction and can deduce that $V=L(\lambda)$ and $L(\lambda)$ is the direct sum of the irreducible $T$-modules $T v$ for $v \in B$, each of which is isomorphic to a Verma module for $T$.

Proposition 20.21 Let $\mu$ be a weight of $L(\lambda)$ where $\lambda \in X^{+}$and $\lambda(c)>0$. Then the multiplicities of the weights $\mu, \mu-\delta$ satisfy $m_{\mu-\delta} \geq m_{\mu}$.

Proof. This follows from Proposition 20.20. We choose a non-zero element $x \in L(A)_{-\delta}$. Consider the action of $x$ on the $T$-module $L(\lambda)$. This $T$-module is a direct sum of Verma modules for $T$. Since $x \in T^{-}, x$ acts on each Verma module for $T$ injectively. Thus $x$ acts on $L(\lambda)$ injectively. We have a map

$$
\begin{aligned}
& L(\lambda)_{\mu} \rightarrow L(\lambda)_{\mu-\delta} \\
& v \rightarrow x v
\end{aligned}
$$

which is injective, and so

$$
\operatorname{dim} L(\lambda)_{\mu-\delta} \geq \operatorname{dim} L(\lambda)_{\mu}
$$

that is $m_{\mu-\delta} \geq m_{\mu}$ as required.
Thus the multiplicities form an increasing sequence as we move down a string of weights for $L(\lambda)$.

### 20.4 The fundamental modules for $L\left(\tilde{A}_{1}\right)$

We now give an example of the situation described in Section 20.3. We consider the affine Kac-Moody algebra of type $\tilde{A}_{1}$. This has diagram

and Cartan matrix

$$
A=\left(\begin{array}{cc}
2 & -2 \\
-2 & 2
\end{array}\right)
$$

We consider the irreducible modules $L\left(\omega_{0}\right), L\left(\omega_{1}\right)$ where $\omega_{0}, \omega_{1}$ are the fundamental weights. By symmetry we need only determine the character of one of these. We shall consider the module $L\left(\omega_{0}\right)$.

In type $\tilde{A}_{1}$ we have $\delta=\alpha_{0}+\alpha_{1}$ and $c=h_{0}+h_{1}$, that is

$$
a_{0}=1, \quad a_{1}=1, \quad c_{0}=1, \quad c_{1}=1
$$

We recall that

$$
H^{*}=\left(H^{0}\right)^{*} \oplus(\mathbb{C} \gamma+\mathbb{C} \delta)
$$

and that

$$
\lambda=\lambda^{0}+\lambda(c) \gamma+a_{0}^{-1} \lambda(d) \delta \quad \text { by Lemma 20.1. }
$$

The root lattice $Q$ is given by

$$
Q=\mathbb{Z} \alpha_{0}+\mathbb{Z} \alpha_{1}
$$

and we have $\alpha_{0}^{0}=-\alpha_{1}, \quad \alpha_{1}^{0}=\alpha_{1}$. We have $\left(H^{0}\right)^{*}=\mathbb{C} \alpha_{1}$ and the lattice $Q^{0} \subset\left(H^{0}\right)^{*}$ is given by $Q^{0}=\mathbb{Z} \alpha_{1}$. We have $\theta=\alpha_{1}$ and $h_{\theta}=\frac{1}{a_{0}}\left(c-h_{0}\right)=h_{1}$. The closure of the fundamental alcove is given by

$$
\begin{aligned}
\bar{A}^{*} & =\left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{1}\right) \geq 0, \lambda\left(h_{\theta}\right) \leq 1\right\} \\
& =\left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; 0 \leq \lambda\left(h_{1}\right) \leq 1\right\}
\end{aligned}
$$

We have $\omega_{0}=\gamma$ and $\gamma^{0}=0$. Thus

$$
\begin{aligned}
\left(\gamma^{0}+Q^{0}\right) \cap \bar{A}^{*} & =\left\{m \alpha_{1} ; m \in \mathbb{Z}, 0 \leq 2 m \leq 1\right\} \\
& =\{0\}
\end{aligned}
$$

Thus by Proposition 20.16 the module $L(\gamma)$ has only one dominant maximal weight, which must be the highest weight $\gamma$. The other maximal weights are the transforms of $\gamma$ under the affine Weyl group $W=\left\langle s_{0}, s_{1}\right\rangle$. We have

$$
\begin{aligned}
& s_{0}\left(\alpha_{0}\right)=-\alpha_{0} \quad s_{0}\left(\alpha_{1}\right)=2 \alpha_{0}+\alpha_{1} \\
& s_{1}\left(\alpha_{0}\right)=\alpha_{0}+2 \alpha_{1} \quad s_{1}\left(\alpha_{1}\right)=-\alpha_{1}
\end{aligned}
$$

The action of $s_{0}, s_{1}$ on the basis $\gamma, \alpha_{1}, \delta$ of $H^{*}$ is given by

$$
\begin{array}{lll}
s_{0}(\gamma)=\gamma+\alpha_{1}-\delta & s_{0}\left(\alpha_{1}\right)=-\alpha_{1}+2 \delta & s_{0}(\delta)=\delta \\
s_{1}(\gamma)=\gamma & s_{1}\left(\alpha_{1}\right)=-\alpha_{1} & s_{1}(\delta)=\delta
\end{array}
$$

The affine Weyl group $W$ is an infinite dihedral group and has a semidirect product decomposition

$$
W=t\left(Q^{0}\right) W^{0}=W^{0} t\left(Q^{0}\right)
$$

where $Q^{0}=M^{*}=\mathbb{Z} \alpha_{1}$ and $W^{0}=\left\{1, s_{1}\right\}$. The translation $t_{\mu}$ for $\mu \in Q^{0}$ is given by

$$
t_{\mu}(\lambda)=\lambda+\lambda(c) \mu-\left(\langle\lambda, \mu\rangle+\frac{1}{2}\langle\mu, \mu\rangle \lambda(c)\right) \delta
$$

which in the present case gives

$$
\begin{aligned}
& t_{m \alpha_{1}}(\gamma)=\gamma+m \alpha_{1}-m^{2} \delta \\
& t_{m \alpha_{1}}\left(\alpha_{1}\right)=\alpha_{1}-2 m \delta \\
& t_{m \alpha_{1}}(\delta)=\delta
\end{aligned}
$$

for $m \in \mathbb{Z}$. The stabiliser of $\gamma$ in $W$ is $W^{0}$ and the maximal weights in $L(\gamma)$ have the form $\gamma+m \alpha_{1}-m^{2} \delta$ for $m \in \mathbb{Z}$. The set of all weights of $L(\gamma)$ is $\gamma+m \alpha_{1}-m^{2} \delta-k \delta$ for $m \in \mathbb{Z}, k \in \mathbb{Z}$ and $k \geq 0$. The weights $\gamma+m \alpha_{1}-m^{2} \delta$


Figure 20.1 Maximal weights in $L(\gamma)$
have multiplicity 1 , and $\gamma+m \alpha_{1}-m^{2} \delta-k \delta$ has multiplicity depending only on $k$ (i.e. independent of $m$ ). The weights are shown in Figure 20.1.

We shall determine the multiplicities of these weights. We use Kac' character formula

$$
\operatorname{ch} L(\gamma)=\frac{\sum_{w \in W} \varepsilon(w) e_{w(\gamma+\rho)-\rho}}{\prod_{\alpha \in \Phi^{+}}\left(1-e_{-\alpha}\right)^{m_{\alpha}}}
$$

Now

$$
\begin{aligned}
\sum_{\omega \in W} \varepsilon(w) e_{w(\gamma+\rho)-\rho} & =\sum_{w^{0} \in W^{0}} \sum_{\mu \in \mathbb{Z} \alpha_{1}} \varepsilon\left(w^{0}\right) e_{w^{0} t_{\mu}(\gamma+\rho)-\rho} \\
& =\sum_{w^{0} \in W^{0}} \varepsilon\left(w^{0}\right) \sum_{n \in \mathbb{Z}} e_{w^{0} t_{n \alpha_{1}}(\gamma+\rho)-\rho}
\end{aligned}
$$

Now $\rho=\rho^{0}+2 \gamma$ by Lemma 20.2 where $\rho^{0}=\frac{1}{2} \alpha_{1}$. Thus $\rho=2 \gamma+\frac{1}{2} \alpha_{1}$ and $\gamma+\rho=3 \gamma+\frac{1}{2} \alpha_{1}$. Hence

$$
t_{n \alpha_{1}}(\gamma+\rho)=3 \gamma+\left(3 n+\frac{1}{2}\right) \alpha_{1}-\left(3 n^{2}+n\right) \delta
$$

so $t_{n \alpha_{1}}(\gamma+\rho)-\rho=\gamma+3 n \alpha_{1}-\left(3 n^{2}+n\right) \delta$. Also

$$
s_{1} t_{n \alpha_{1}}(\gamma+\rho)=3 \gamma-\left(3 n+\frac{1}{2}\right) \alpha_{1}-\left(3 n^{2}+n\right) \delta
$$

so $s_{1} t_{n \alpha_{1}}(\gamma+\rho)-\rho=\gamma-(3 n+1) \alpha_{1}-\left(3 n^{2}+n\right) \delta$. Thus

$$
\sum_{w \in W} \varepsilon(w) e_{w(\gamma+\rho)-\rho}=e_{\gamma} \sum_{n \in \mathbb{Z}}\left(e_{3 n \alpha_{1}}-e_{-(3 n+1) \alpha_{1}}\right) e_{-\left(3 n^{2}+n\right) \delta}
$$

We write $e_{-\alpha_{1}}=z$ and $e_{-\delta}=q^{1 / 2}$. Then our expression is

$$
e_{\gamma} \sum_{n \in \mathbb{Z}}\left(z^{-3 n}-z^{3 n+1}\right) q^{n(3 n+1) / 2}
$$

Now we may factorise this expression by using Macdonald's identity for type $\tilde{A}_{1}^{\prime}$. By Proposition 20.6 it is equal to

$$
\begin{aligned}
& e_{\gamma} \prod_{n>0}\left(1-q^{n}\right)\left(1-q^{n} z^{-1}\right)\left(1-q^{n-1} z\right)\left(1-q^{2 n-1} z^{-2}\right)\left(1-q^{2 n-1} z^{2}\right) \\
& =e_{\gamma}(1-z) \prod_{n>0}\left(1-q^{n}\right)\left(1-q^{n} z^{-1}\right)\left(1-q^{n} z\right)\left(1-q^{2 n-1} z^{-2}\right)\left(1-q^{2 n-1} z^{2}\right) \\
& =e_{\gamma}(1-z) \prod_{n>0}\left(1-q^{n}\right)\left(1-q^{n} z^{-1}\right)\left(1-q^{\frac{2 n-1}{2}} z^{-1}\right)\left(1-q^{n} z\right)\left(1-q^{\frac{2 n-1}{2}} z\right) \\
& \quad \times\left(1+q^{\frac{2 n-1}{2}} z^{-1}\right)\left(1+q^{\frac{2 n-1}{2}} z\right) \\
& =e_{\gamma}(1-z) \prod_{k>0}\left(1-q^{k / 2} z^{-1}\right)\left(1-q^{k / 2} z\right) \prod_{n>0}\left(1-q^{n}\right) \\
& \quad \times\left(1+q^{\frac{2 n-1}{2}} z^{-1}\right)\left(1+q^{\frac{2 n-1}{2}} z\right) .
\end{aligned}
$$

We now make use of Macdonald's identity for type $\tilde{A}_{1}$. By Proposition 20.5 this asserts that

$$
\prod_{n>0}\left(1-q^{n}\right)\left(1-q^{n-1} z^{\prime}\right)\left(1-q^{n} z^{\prime-1}\right)=\sum_{n>0}(-1)^{n}\left(z^{\prime n}-z^{\prime-(n-1)}\right) q^{\frac{n(n-1)}{2}}
$$

Putting $z^{\prime}=-z^{-1} q^{\frac{1}{2}}$ we obtain

$$
\begin{aligned}
\prod_{n>0}\left(1-q^{n}\right)\left(1+q^{\frac{2 n-1}{2}} z^{-1}\right)\left(1+q^{\frac{2 n-1}{2}} z\right) & =\sum_{n>0}\left(z^{-n} q^{n^{2} / 2}+z^{n-1} q^{(n-1)^{2} / 2}\right) \\
& =\sum_{n \in \mathbb{Z}} z^{-n} q^{n^{2} / 2}
\end{aligned}
$$

Hence

$$
\begin{aligned}
\operatorname{ch} L(\gamma) & =\frac{\sum_{w \in W} \varepsilon(w) e_{w(\gamma+\rho)-\rho}}{(1-z) \prod_{k>0}\left(1-q^{k / 2} z^{-1}\right)\left(1-q^{k / 2} z\right)\left(1-q^{k / 2}\right)} \\
& =\frac{e_{\gamma} \sum_{n \in \mathbb{Z}} z^{-n} q^{n^{2} / 2}}{\prod_{k>0}\left(1-q^{k / 2}\right)}=\frac{\sum_{n \in \mathbb{Z}} e_{\gamma+n \alpha_{1}-n^{2} \delta}}{\prod_{k>0}\left(1-e_{-k \delta}\right)}
\end{aligned}
$$

Now

$$
\begin{aligned}
\frac{1}{\prod_{k>0}\left(1-e_{-k \delta}\right)} & =\prod_{k>0}\left(1+e_{-k \delta}+e_{-2 k \delta}+\cdots\right) \\
& =\sum_{k \geq 0} p(k) e_{-k \delta}
\end{aligned}
$$

where $p(k)$ is the number of partitions of $k$. Thus

$$
\operatorname{ch} L(\gamma)=\sum_{n \in \mathbb{Z}} \sum_{k \geq 0} p(k) e_{\gamma+n \alpha_{1}-n^{2} \delta-k \delta}
$$

Hence we have proved
Proposition 20.22 The weights of the fundamental module $L(\gamma)$ for $L\left(\tilde{A}_{1}\right)$ are $\gamma+n \alpha_{1}-n^{2} \delta-k \delta$ for $n \in \mathbb{Z}$ and $k \geq 0$. This weight has multiplicity $p(k)$.

We note in particular that all the maximal weights $\gamma+n \alpha_{1}-n^{2} \delta$ have multiplicity 1 and that the multiplicity of the weight $\mu-k \delta$ in the string with maximal weight $\mu$ depends only upon $k$ and not on $\mu$.

### 20.5 The basic representation

The module $L\left(\omega_{0}\right)$ for an affine Kac-Moody algebra $L(A)$ gives the so-called basic representation of $L(A)$. Since $\omega_{0}=\gamma$ we have described the character of the basic representation of $L\left(\tilde{A}_{1}\right)$. We shall state without proof some generalisations of this character formula to other types of affine Kac-Moody algebras. For simplicity we shall concentrate on those of types $\tilde{A}_{l}, \tilde{D}_{l}$ and $\tilde{E}_{l}$.

Theorem 20.23 The basic representation $L(\gamma)$ for the Kac-Moody algebra $L(A)$ of types $\tilde{A}_{l}, \tilde{D}_{l}, \tilde{E}_{6}, \tilde{E}_{7}, \tilde{E}_{8}$ has the following properties:
(a) $\gamma$ is the unique dominant maximal weight of $L(\gamma)$.
(b) The set of all maximal weights is

$$
\left\{\gamma+\mu-\frac{1}{2}\langle\mu, \mu\rangle \delta \quad \text { for } \mu \in Q^{0}\right\} .
$$

(c) The set of all weights is

$$
\left\{\gamma+\mu-\frac{1}{2}\langle\mu, \mu\rangle \delta-k \delta \quad \text { for } \mu \in Q^{0}, \quad k \in \mathbb{Z}, \quad k \geq 0\right\} .
$$

(d) The character of the basic representation is

$$
\operatorname{ch} L(\gamma)=\frac{\sum_{\mu \in Q^{0}} e_{\gamma+\mu-\frac{1}{2}\langle\mu, \mu\rangle \delta}}{\left(\prod_{k>0}\left(1-q^{k}\right)\right)^{l}}
$$

where $q=e^{-\delta}$.
(e) The multiplicity of the weight $\gamma+\mu-\frac{1}{2}\langle\mu, \mu\rangle \delta-k \delta$ is $p_{l}(k)$, the number of partitions of $k$ into $l$ colours. We have

$$
\frac{1}{\left(\prod_{k>0}\left(1-q^{k}\right)\right)^{l}}=\sum_{k \geq 0} p_{l}(k) q^{k} .
$$

The proof of this theorem can be found in the book of Kac, InfiniteDimensional Lie Algebras, third edition, Chapter 12.

We shall also describe without proof how to obtain a realisation of the basic representation $L(\gamma)$ of $L(A)$ in types $\tilde{A}_{l}, \tilde{D}_{l}, \tilde{E}_{l}$. We first make some comments on differential operators. Let $R=\mathbb{C}\left[x_{1}, x_{2}, x_{3}, \ldots\right]$ be the polynomial ring over $\mathbb{C}$ in countably many variables and $\hat{R}=\mathbb{C}\left[\left[x_{1}, x_{2}, x_{3}, \ldots\right]\right]$ be the ring of formal power series in these variables. We shall consider differential operators on $R$ with values in $\hat{R}$. An example is the partial derivative $\partial / \partial x_{i}$ or, more generally, the divided power $\frac{1}{m_{i}!}\left(\partial / \partial x_{i}\right)^{m_{i}}$. We also have finite products $\prod_{i} \frac{1}{m_{i}!}\left(\partial / \partial x_{i}\right)^{m_{i}}$ where $m=\left(m_{1}, m_{2}, m_{3}, \ldots\right)$ satisfies the conditions that $m_{i} \in$ $\mathbb{Z}, m_{i} \geq 0$, and $m_{i}>0$ for only finitely many $i$. We define $D_{m}: R \rightarrow \hat{R}$ by

$$
D_{m}=\prod_{i} \frac{1}{m_{i}!}\left(\partial / \partial x_{i}\right)^{m_{i}} .
$$

We also allow such operators combined with multiplication by elements of $\hat{R}$. Thus

$$
\sum_{m} P_{m} D_{m}: R \rightarrow \hat{R}
$$

is a differential operator, where $P_{m} \in \hat{R}$ and the sum over $\boldsymbol{m}$ will in general be infinite. $\sum_{m} P_{m} D_{m}$ is a linear map from $R$ to $\hat{R}$. In fact each linear map from $R$ to $\hat{R}$ has this form, as we now show.

Proposition 20.24 Each linear map from $R$ to $\hat{R}$ can be written as $\sum_{m} P_{m} D_{m}$ for a unique set of elements $P_{m} \in \hat{R}$.

Proof. Let $M_{m} \in R$ be the monomial $M_{m}=\prod_{i} x_{i}^{m_{i}}$. The monomials $M_{m}$ form a basis for $R$. We have $D_{m}\left(M_{k}\right)=0$ unless $k_{i} \geq m_{i}$ for each $i$. We write this condition as $\boldsymbol{k} \geq \boldsymbol{m}$. We write $\boldsymbol{k}>\boldsymbol{m}$ if $\boldsymbol{k} \geq \boldsymbol{m}$ and $\boldsymbol{k} \neq \boldsymbol{m}$. We also have

$$
D_{\boldsymbol{m}}\left(M_{\boldsymbol{k}}\right)=\binom{\boldsymbol{k}}{\boldsymbol{m}} M_{\boldsymbol{k}-\boldsymbol{m}} \quad \text { if } \boldsymbol{k} \geq \boldsymbol{m}
$$

where $\binom{k}{m}=\prod_{i}\binom{k_{i}}{m_{i}}$ and $\binom{0}{0}=1$.
Let $\Delta: R \rightarrow \hat{R}$ be the linear map given by $\Delta\left(M_{m}\right)=Q_{m} \in \hat{R}$. We show $\Delta$ is uniquely expressible in the form $\sum P_{m} D_{m}$. We have

$$
\begin{aligned}
\sum_{m} P_{m} D_{m}\left(M_{k}\right) & =\sum_{m \leq k} P_{m}\binom{\boldsymbol{k}}{\boldsymbol{m}} M_{k-m} \\
& =P_{k}+\sum_{m<k} P_{m}\binom{\boldsymbol{k}}{\boldsymbol{m}} M_{k-m} .
\end{aligned}
$$

The condition we require on the $P_{m}$ is that

$$
P_{k}+\sum_{m<k} P_{m}\binom{\boldsymbol{k}}{\boldsymbol{m}} M_{k-m}=Q_{k} \quad \text { for all } \boldsymbol{k}
$$

In particular $P_{0}=Q_{0}$. Assuming inductively that $P_{m}$ is uniquely determined for all $\boldsymbol{m}<\boldsymbol{k}$ we conclude that

$$
P_{k}=Q_{k}-\sum_{m<k} P_{m}\binom{\boldsymbol{k}}{\boldsymbol{m}} M_{k-m}
$$

is uniquely determined.
Thus the set of differential operators from $R$ to $\hat{R}$ is the set of all linear maps from $R$ to $\hat{R}$.

We shall now consider certain special kinds of differential operators. Let $\lambda=\left(\lambda_{1}, \lambda_{2}, \lambda_{3}, \ldots\right)$ where $\lambda_{i} \in \mathbb{C}$. Here there may be infinitely many non-zero $\lambda_{i}$. Define $T_{\lambda}: R \rightarrow \hat{R}$ by

$$
T_{\lambda} f\left(x_{1}, x_{2}, x_{3}, \ldots\right)=f\left(x_{1}+\lambda_{1}, x_{2}+\lambda_{2}, x_{3}+\lambda_{3}, \ldots\right)
$$

$T_{\lambda}$ is clearly a linear map from $R$ to $\hat{R}$. It may be written as a differential operator by using the Taylor expansion. We have

$$
\begin{aligned}
f\left(x_{1}\right. & \left.+\lambda_{1}, x_{2}+\lambda_{2}, x_{3}+\lambda_{3}, \ldots\right) \\
& =\sum_{\boldsymbol{m}} \frac{\lambda_{1}^{m_{1}}}{m_{1}!} \frac{\lambda_{2}^{m_{2}}}{m_{2}!} \frac{\lambda_{3}^{m_{3}}}{m_{3}!} \ldots\left(\partial / \partial x_{1}\right)^{m_{1}}\left(\partial / \partial x_{2}\right)^{m_{2}}\left(\partial / \partial x_{3}\right)^{m_{3}} \ldots f\left(x_{1}, x_{2}, x_{3}, \ldots\right) \\
& =\sum_{\boldsymbol{m}}\left(\prod_{i} \lambda_{i}^{m_{i}}\right) D_{m} f\left(x_{1}, x_{2}, x_{3}, \ldots\right) .
\end{aligned}
$$

Thus $T_{\lambda}=\sum_{m}\left(\prod_{i} \lambda_{i}^{m_{i}}\right) D_{m}$. The operator $T_{\lambda}$ may also be written in the following convenient form. We have

$$
\begin{aligned}
f\left(x_{1}\right. & \left.+\lambda_{1}, x_{2}+\lambda_{2}, x_{3}+\lambda_{3}, \ldots\right) \\
& =\left(\sum_{m_{1}} \frac{\lambda_{1}^{m_{1}}}{m_{1}!}\left(\partial / \partial x_{1}\right)^{m_{1}}\right)\left(\sum_{m_{2}} \frac{\lambda_{2}^{m_{2}}}{m_{2}!}\left(\partial / \partial x_{2}\right)^{m_{2}}\right) \ldots f\left(x_{1}, x_{2}, x_{3}, \ldots\right) \\
& =\exp \left(\lambda_{1} \partial / \partial x_{1}\right) \exp \left(\lambda_{2} \partial / \partial x_{2}\right) \ldots f\left(x_{1}, x_{2}, x_{3}, \ldots\right) \\
& =\exp \left(\lambda_{1} \partial / \partial x_{1}+\lambda_{2} \partial / \partial x_{2}+\ldots\right) f\left(x_{1}, x_{2}, x_{3}, \ldots\right)
\end{aligned}
$$

Thus $T_{\lambda}=\exp \left(\sum_{i} \lambda_{i}\left(\partial / \partial x_{i}\right)\right)$.

Lemma 20.25 Suppose $D: R \rightarrow \hat{R}$ is a linear map which satisfies $\left[x_{i}, D\right]=$ $\lambda_{i} D$ for all $i$, that is

$$
x_{i} D f-D\left(x_{i} f\right)=\lambda_{i} D f \quad \text { for all } f \in R
$$

Then

$$
D=D(1) \exp \left(-\sum_{i} \lambda_{i} \frac{\partial}{\partial x_{i}}\right)
$$

Proof. We shall show $D f=D(1) \exp \left(-\sum_{i} \lambda_{i}\left(\partial / \partial x_{i}\right)\right) f$ for all monomials $f \in R$, using induction on the degree of $f$. If $f$ has degree 0 then $f=c \in \mathbb{C}$ and we have

$$
D(1) \exp \left(-\sum_{i} \lambda_{i} \frac{\partial}{\partial x_{i}}\right) c=D(1) c=D(c)
$$

Assuming the result for a monomial $f$ we prove it for $x_{i} f$. We have

$$
D\left(x_{i} f\right)=\left(x_{i}-\lambda_{i}\right) D f=\left(x_{i}-\lambda_{i}\right) D(1) \exp \left(-\sum_{i} \lambda_{i} \frac{\partial}{\partial x_{i}}\right) f
$$

On the other hand

$$
\begin{aligned}
D(1) \exp \left(-\sum_{i} \lambda_{i} \frac{\partial}{\partial x_{i}}\right)\left(x_{i} f\right) & =D(1) T_{-\lambda}\left(x_{i} f\right)=D(1)\left(x_{i}-\lambda_{i}\right) T_{-\lambda} f \\
& =D(1)\left(x_{i}-\lambda_{i}\right) \exp \left(-\sum_{i} \lambda_{i} \frac{\partial}{\partial x_{i}}\right) f
\end{aligned}
$$

Thus the lemma is proved.

Lemma 20.26 Suppose $D: R \rightarrow \hat{R}$ is a linear map which satisfies $\left[\partial / \partial x_{i}, D\right]=$ $\mu_{i} D$ for all $i$, that is

$$
\left(\partial / \partial x_{i}\right)(D f)-D\left(\partial f / \partial x_{i}\right)=\mu_{i} D f \quad \text { for all } f \in R
$$

Then

$$
D(1)=c \exp \left(\sum_{i} \mu_{i} x_{i}\right) \quad \text { for some } c \in \mathbb{C} .
$$

Proof. Consider the element $\exp \left(\sum_{i}-\mu_{i} x_{i}\right) D f \in \hat{R}$. We have

$$
\begin{aligned}
& \frac{\partial}{\partial x_{i}}\left(\exp \left(\sum_{i}-\mu_{i} x_{i}\right) D f\right) \\
& \quad=-\mu_{i} \exp \left(\sum_{i}-\mu_{i} x_{i}\right) D f+\exp \left(\sum_{i}-\mu_{i} x_{i}\right) \frac{\partial}{\partial x_{i}}(D f) \\
& \quad=\exp \left(\sum_{i}-\mu_{i} x_{i}\right)\left(\partial / \partial x_{i}-\mu_{i}\right) D f
\end{aligned}
$$

Thus by the assumption of the lemma we have

$$
\begin{aligned}
\left(\exp \left(\sum_{i}-\mu_{i} x_{i}\right) D\right) \partial f / \partial x_{i} & =\exp \left(\sum_{i}-\mu_{i} x_{i}\right)\left(\partial / \partial x_{i}-\mu_{i}\right) D f \\
& =\frac{\partial}{\partial x_{i}}\left(\exp \left(\sum_{i}-\mu_{i} x_{i}\right) D f\right)
\end{aligned}
$$

Write $\Delta=\exp \left(\sum_{i}-\mu_{i} x_{i}\right) D$. Then we have $\Delta \partial f / \partial x_{i}=\left(\partial / \partial x_{i}\right)(\Delta f)$ for each $i$ and $f$. In particular we may put $f=1$ and obtain $\left(\partial / \partial x_{i}\right)(\Delta(1))=0$. Thus $\Delta(1)=c$ for some $c \in \mathbb{C}$. Hence

$$
D(1)=\exp \left(\sum_{i} \mu_{i} x_{i}\right) \Delta(1)=c \exp \left(\sum_{i} \mu_{i} x_{i}\right)
$$

as required.
Proposition 20.27 The set of all differential operators $D: R \rightarrow \hat{R}$ satisfying the conditions $\left[x_{i}, D\right]=\lambda_{i} D$ and $\left[\partial / \partial x_{i}, D\right]=\mu_{i} D$ for $\lambda_{i}, \mu_{i} \in \mathbb{C}$ forms a 1-dimensional vector space with basis

$$
\exp \left(\sum \mu_{i} x_{i}\right) \exp \left(-\sum \lambda_{i} \frac{\partial}{\partial x_{i}}\right)
$$

Proof. This follows from Lemmas 20.25 and 20.26.
Definition Differential operators $D: R \rightarrow \hat{R}$ of the form

$$
\exp \left(\sum \mu_{i} x_{i}\right) \exp \left(-\sum \lambda_{i} \frac{\partial}{\partial x_{i}}\right)
$$

for $\lambda_{i}, \mu_{i} \in \mathbb{C}$ are called vertex operators.
Now let $L=\hat{\mathfrak{R}}\left(L^{0}\right)$ be an affine Kac-Moody algebra of type $\tilde{A}_{l}, \tilde{D}_{l}$ or $\tilde{E}_{l}$. Let

$$
T^{-}=\bigoplus_{j<0} L_{j \delta}
$$

Then $T^{-}$has a basis $t^{j} \otimes h_{i}$ for $i=1, \ldots l$ and $j<0$. Consider the symmetric algebra $S\left(T^{-}\right)$. This is isomorphic to the polynomial ring over $\mathbb{C}$ in the variables $t^{j} \otimes h_{i}$.

Let $Q^{0}$ be the subgroup of $H^{0}$ generated by $h_{1}, \ldots, h_{l}$. We shall write $Q^{0}$ multiplicatively, so that its elements have form $h_{1}^{m_{1}} \ldots h_{l}^{m_{l}}$ with $m_{1}, \ldots, m_{l} \in$ $\mathbb{Z}$. Let $\mathbb{C}\left[Q^{0}\right]$ be the group algebra of $Q^{0}$ over $\mathbb{C}$. Elements of $\mathbb{C}\left[Q^{0}\right]$ have form

$$
\sum_{m_{1}, \ldots, m_{l} \in \mathbb{Z}} \lambda_{m_{1}, \ldots, m_{l}} h_{1}^{m_{1}} \ldots h_{l}^{m_{l}}
$$

$\mathbb{C}\left[Q^{0}\right]$ is isomorphic to the algebra of Laurent polynomials over $\mathbb{C}$ in the variables $h_{1}, \ldots, h_{l}$.

We now form the tensor product

$$
V=S\left(T^{-}\right) \otimes \mathbb{C}\left[Q^{0}\right]
$$

$V$ is isomorphic to the algebra

$$
\mathbb{C}\left[h_{1}, \ldots, h_{l}, h_{1}^{-1}, \ldots, h_{l}^{-1}, \quad t^{j} \otimes h_{i}\right]
$$

for $i=1, \ldots, l$ and $j<0$.
We define certain maps $h_{\alpha}(n): V \rightarrow V$ out of which vertex operators will be constructed. For $n \in \mathbb{Z}$ with $n>0, h_{\alpha}(n)$ is the derivation of $V$ uniquely determined by the conditions

$$
\begin{aligned}
& t^{-n} \otimes h_{i} \rightarrow n\left\langle h_{i}, h_{\alpha}\right\rangle \\
& t^{j} \otimes h_{i} \rightarrow 0 \quad \text { for } j \neq-n \\
& h_{i} \rightarrow 0
\end{aligned}
$$

For $n \in \mathbb{Z}$ with $n<0, h_{\alpha}(n): V \rightarrow V$ is multiplication by $\left(t^{n} \otimes h_{\alpha}\right) \otimes 1$.
We now consider the expression

$$
\exp \left(\sum_{n<0}-\frac{h_{\alpha}(n)}{n} z^{-n}\right) \exp \left(\sum_{n>0}-\frac{h_{\alpha}(n)}{n} z^{-n}\right)
$$

where $z$ is an indeterminate. We first observe that $\exp \left(\sum_{n>0}-\frac{h_{\alpha}(n)}{n} z^{-n}\right)$ maps $V$ into $\mathbb{C}\left[z^{-1}\right] \otimes V$. To see this we observe that each element $v \in V$ is a finite linear combination of monomials

$$
M_{\boldsymbol{m}}=\prod_{\substack{i, j \\ j<0}}\left(t^{j} \otimes h_{i}\right)^{m_{i j}} \prod_{i} h_{i}^{m_{i}}
$$

where $\boldsymbol{m}=\left(m_{i j}, m_{i}\right)$ with $m_{i} \in \mathbb{Z}, m_{i j} \in \mathbb{Z}, m_{i j} \geq 0$.

Let $d_{1}(\boldsymbol{m})=\sum_{i, j} m_{i j}$. Then the derivation $h_{\alpha}(n), n>0$, transforms $M_{m}$ into a linear combination of monomials in which $d_{1}(\boldsymbol{m})$ is decreased by 1 and $\prod_{i} h_{i}^{m_{i}}$ remains unchanged. Thus a succession of $d_{1}(\boldsymbol{m})+1$ derivations $h_{\alpha}(n)$ for various $n>0$ annihilates $M_{m}$. Also, for a given monomial $M_{m}, h_{\alpha}(n)$ annihilates $M_{m}$ for all but finitely many $n>0$. Thus in the expression

$$
\exp \left(\sum_{n>0}-\frac{h_{\alpha}(n)}{n} z^{-n}\right) v \quad v \in V
$$

only finitely many terms $-\frac{h_{\alpha}(n)}{n} z^{-n}$ act on $v$ and only a finite set of products of such terms can act on $v$ to give a non-zero element. Thus we have

$$
\exp \left(\sum_{n>0}-\frac{h_{\alpha}(n)}{n} z^{-n}\right): V \rightarrow \mathbb{C}\left[z^{-1}\right] \otimes V
$$

We shall modify this operator in the following way. Let

$$
\varepsilon: Q^{0} \times Q^{0} \rightarrow\{ \pm 1\}
$$

be the function defined by

$$
\begin{aligned}
& \varepsilon\left(h_{i}, h_{i}\right)=-1 \\
& \varepsilon\left(h_{i}, h_{j}\right)=1 \quad \text { if } A_{i j}=0 \\
& \varepsilon\left(h_{i}, h_{j}\right) \text { is given by }\left[E_{\alpha_{i}} E_{\alpha_{j}}\right]=\varepsilon\left(h_{i}, h_{j}\right) E_{\alpha_{i}+\alpha_{j}} \text { if } A_{i j}=-1 \\
& \varepsilon\left(h+h^{\prime}, h^{\prime \prime}\right)=\varepsilon\left(h, h^{\prime \prime}\right) \varepsilon\left(h^{\prime}, h^{\prime \prime}\right) \\
& \varepsilon\left(h, h^{\prime}+h^{\prime \prime}\right)=\varepsilon\left(h, h^{\prime}\right) \varepsilon\left(h, h^{\prime \prime}\right)
\end{aligned}
$$

Given $\alpha \in \Phi^{0}$ we define a map $\varepsilon^{\alpha} \in$ End $V$ by

$$
\varepsilon^{\alpha}: P \otimes h \rightarrow P \otimes \varepsilon\left(h_{\alpha}, h\right) h
$$

where $P \in S\left(T^{-}\right), h \in Q^{0}$. We also define $e^{\alpha} \in$ End $V$ by

$$
e^{\alpha}: P \otimes h \rightarrow P \otimes h_{\alpha} h
$$

and $z^{\alpha} \in \operatorname{End}\left(\mathbb{C}\left[z, z^{-1}\right] \otimes V\right)$ by

$$
z^{\alpha}:(P \otimes h) z^{i} \rightarrow(P \otimes h) z^{i+\left\langle h_{\alpha}, h\right\rangle}
$$

We now define, for $\alpha \in \Phi^{0}$, the operator $Y_{\alpha}(z)$ on $V$ by

$$
Y_{\alpha}(z)=\exp \left(\sum_{n<0}-\frac{h_{\alpha}(n)}{n} z^{-n}\right) \exp \left(\sum_{n>0}-\frac{h_{\alpha}(n)}{n} z^{-n}\right) e^{\alpha} z^{\alpha} \varepsilon^{\alpha}
$$

We claim that $Y_{\alpha}(z)$ can be written in the form

$$
Y_{\alpha}(z)=\sum_{j \in \mathbb{Z}} \Gamma_{\alpha}(j) z^{-j-1}
$$

where $\Gamma_{\alpha}(j) \in$ End $V$. In order to see this we consider the effect of $Y_{\alpha}(z)$ on a monomial $M_{m}$ in $V$. We have

$$
e^{\alpha} z^{\alpha} \varepsilon^{\alpha} M_{\boldsymbol{m}} \in z^{n_{\alpha}} \otimes V
$$

where $n_{\alpha}=\sum_{i} m_{i}\left\langle h_{\alpha}, h_{i}\right\rangle$. Thus

$$
\exp \left(\sum_{n>0}-\frac{h_{\alpha}(n)}{n} z^{-n}\right) e^{\alpha} z^{\alpha} \varepsilon^{\alpha} M_{m} \in \sum_{k=0}^{K} z^{n_{\alpha}-k} \otimes V
$$

for some $K>0$, and

$$
\begin{aligned}
& \exp \left(\sum_{n<0}-\frac{h_{\alpha}(n)}{n} z^{-n}\right) \exp \left(\sum_{n>0}-\frac{h_{\alpha}(n)}{n} z^{-n}\right) e^{\alpha} z^{\alpha} \varepsilon^{\alpha} M_{m} \\
& \quad \in \sum_{k^{\prime} \geq 0} \sum_{k=0}^{K} z^{n_{\alpha}-k+k^{\prime}} \otimes V
\end{aligned}
$$

Thus to obtain $z^{-j-1}$ on the right-hand side of $Y_{\alpha}(z) M_{m}$ we must have $n_{\alpha}-k+k^{\prime}=-j-1$. For each value of $k$ there is at most one $k^{\prime} \geq 0$ satisfying this. Since only finitely many $k$ arise, only finitely many $k^{\prime}$ can arise for given $j$. This shows that only finitely many terms $-\frac{h_{\alpha}(n)}{n} z^{-n}$ in $\exp \left(\sum_{n<0}-\frac{h_{\alpha}(n)}{n} z^{-n}\right)$ are involved in $\Gamma_{\alpha}(j)$ and only finitely many products of such terms are involved. Thus we have

$$
\Gamma_{\alpha}(j): V \rightarrow V
$$

and

$$
Y_{\alpha}(z)=\sum_{j \in \mathbb{Z}} \Gamma_{\alpha}(j) z^{-j-1}
$$

with $\Gamma_{\alpha}(j) \in$ End $V$.
Now the vector space $V$ can be regarded as an $L$-module giving the basic representation of $L$. In order to describe the $L$-action on $V$ we introduce some further notation. We have defined $h_{\alpha}(n) \in$ End $V$ for $n>0$ and $n<0$. We now define $h(0) \in$ End $V$ for any $h \in H^{0}$. In contrast to the $h_{\alpha}(n)$ for $n \neq 0$, which act non-trivially on $S\left(T^{-}\right)$and trivially on $\mathbb{C}\left[Q^{0}\right], h(0)$ acts trivially on $S\left(T^{-}\right)$and non-trivially on $\mathbb{C}\left[Q^{0}\right]$. We define, for $h \in H^{0}, h(0): V \rightarrow V$ by

$$
h(0): P \otimes h_{\alpha} \rightarrow P \otimes\left\langle h_{\alpha}, h\right\rangle h_{\alpha}
$$

for $P \in S\left(T^{-}\right), \alpha \in Q^{0}$.

Let $h_{1}^{\prime}, \ldots, h_{l}^{\prime}$ be a basis for $H^{0}$ and $h_{1}^{\prime \prime}, \ldots, h_{l}^{\prime \prime}$ be the dual basis satisfying $\left\langle h_{i}^{\prime}, h_{j}^{\prime \prime}\right\rangle=\delta_{i j}$. We define $D_{0} \in$ End $V$ by

$$
D_{0}=\sum_{i=1}^{l} \frac{1}{2} h_{i}^{\prime}(0) h_{i}^{\prime \prime}(0)+\sum_{n \geq 1} h_{i}^{\prime}(-n) h_{i}^{\prime \prime}(n)
$$

since, for $v \in V, h_{i}^{\prime \prime}(n) v=0$ for all but finitely many $n>0 . D_{0}$ lies in End $V$ and is readily seen to be independent of the choice of basis of $H^{0}$.

We now have the definitions necessary to describe the action of $L$ on $V$ which gives the basic representation.

Theorem 20.28 The vector space $V=S\left(T^{-}\right) \otimes \mathbb{C}\left[Q^{0}\right]$ is a module for the Kac-Moody algebra $L(A)$ of type $\tilde{A}_{l}, \tilde{D}_{l}$ or $\tilde{E}_{l}$ giving the basic representation under the following action $L(A) \rightarrow$ End $V$ :

$$
\begin{aligned}
& t^{n} \otimes H_{\alpha} \rightarrow H_{\alpha}(n) \quad \text { for } \alpha \in \Phi^{0}, n \in \mathbb{Z} \\
& t^{n} \otimes E_{\alpha} \rightarrow \Gamma_{\alpha}(n) \quad \text { for } \alpha \in \Phi^{0}, n \in \mathbb{Z} \\
& c \rightarrow 1_{V} \\
& d \rightarrow-D_{0}
\end{aligned}
$$

The proof of this result can be found in the book of Kac, Infinite-Dimensional Lie Algebras, third edition, Chapter 14.

The highest weight vector of $V$ is the element $1 \otimes 1$. The first 1 is the unit element of the symmetric algebra $S\left(T^{-}\right)$and the second 1 is the unit element of the lattice $Q^{0}$ written multiplicatively, which is the unit element of the group algebra $\mathbb{C}\left[Q^{0}\right]$. This vector $1 \otimes 1$ is annihilated by the generators $e_{i}, i=0,1, \ldots, l$ of $L(A)$. To see this we recall that

$$
e_{i}=1 \otimes E_{i} \quad i=1, \ldots, l \quad e_{0}=t \otimes E_{0}
$$

with $E_{0} \in L_{-\theta}^{0}$. We have

$$
e_{i}(1 \otimes 1)=\Gamma_{\alpha_{i}}(0)(1 \otimes 1) \quad i=1, \ldots, l
$$

which is the coefficient of $z^{-1}$ in $Y_{\alpha_{i}}(z)(1 \otimes 1)$. Also

$$
e_{0}(1 \otimes 1)=\Gamma_{-\theta}(1)(1 \otimes 1)
$$

which is the coefficient of $z^{-2}$ in $Y_{-\theta}(z)(1 \otimes 1)$. Recalling that

$$
\gamma_{\alpha}(z)=\exp \left(\sum_{n<0} \frac{-h_{\alpha}(n)}{n} z^{-n}\right) \exp \left(\sum_{n>0} \frac{-h_{\alpha}(n)}{n} z^{-n}\right) e^{\alpha} z^{\alpha} \varepsilon^{\alpha}
$$

we first note that $e^{\alpha} z^{\alpha} \varepsilon^{\alpha}(1 \otimes 1)=1 \otimes h_{\alpha}$. Now negative powers of $z$ in $Y_{\alpha}(z)(1 \otimes 1)$ can only arise from derivations $h_{\alpha}(n)$ with $n>0$. However, $h_{\alpha}(n)\left(1 \otimes h_{\alpha}\right)=0$ for all $n>0$ since any derivation annihilates the unit element $1 \in S\left(T^{-}\right)$. Thus we have

$$
\begin{aligned}
\Gamma_{\alpha_{i}}(0)(1 \otimes 1) & =0 \quad \text { for } i=1, \ldots, l \\
\Gamma_{-\theta}(1)(1 \otimes 1) & =0
\end{aligned}
$$

Hence $e_{i}(1 \otimes 1)=0$ for all $i=0,1, \ldots, l$.
We now check how the elements $h_{0}, h_{1}, \ldots, h_{l} \in H$ act on $1 \otimes 1$. We have $h_{i}=1 \otimes H_{i}$ for $i=1, \ldots, l$. Thus

$$
h_{i}(1 \otimes 1)=H_{\alpha_{i}}(0)(1 \otimes 1)=1 \otimes\left\langle 0, h_{i}\right\rangle 1=0 .
$$

(We note that in the scalar product $\langle$,$\rangle the elements of Q^{0}$ are written additively so that the unit element will be 0 .) We also have

$$
c=h_{0}+c_{1} h_{1}+\cdots+c_{l} h_{l}
$$

Thus

$$
h_{0}(1 \otimes 1)=c(1 \otimes 1)=1 \otimes 1
$$

Hence we have

$$
h_{i}(1 \otimes 1)=\gamma\left(h_{i}\right)(1 \otimes 1) \quad \text { for } i=0,1, \ldots, l
$$

since $\gamma\left(h_{i}\right)=0$ for $i=1, \ldots, l$ and $\gamma\left(h_{0}\right)=1$. The highest weight vector $v_{\gamma}=1 \otimes 1$ is often called the vacuum vector of the basic representation.

The realisation of the basic representation given by the module $S\left(T^{-}\right) \otimes$ $\mathbb{C}\left[Q^{0}\right]$ is called the homogeneous realisation. It is one of a number of descriptions of the basic representation.

The basic representation is of great importance in a number of applications of the theory of affine Kac-Moody algebras in mathematics and physics. For example applications to the theory of differential equations are described in Kac' book, Chapter 14. There are also particularly interesting applications in the area of mathematical physics. Vertex operators arose in the context of dual resonance models, which subsequently developed into string theory, and the representation theory of affine Kac-Moody algebras plays a key role in string theory. This involves the calculus of vertex operators. The theory of modular forms also plays a key role.

Readers wishing to learn more about the relations between Kac-Moody algebras and string theory may wish to study the 30-page introduction to the book of Frenkel, Lepowsky and Meurman, Vertex Operator Algebras and the

Monster which also explains the connections with modular forms and sporadic simple groups such as the Monster. The book by Kac on Vertex Algebras for Beginners is a useful introduction to the calculus of vertex operators. This whole area relating mathematics and physics is of great current interest and seems certain to continue its rapid development.

## 21

## Borcherds Lie algebras

### 21.1 Definition and examples of Borcherds algebras

A theory of generalised Kac-Moody algebras was introduced by R. Borcherds in 1988. The purpose for which these algebras were introduced was as part of Borcherds' proof of the Conway-Norton conjectures on the representation theory of the Monster simple group, for which Borcherds was awarded a Fields Medal in 1998. These generalised Kac-Moody algebras are now frequently called Borcherds algebras. A detailed discussion of Borcherds algebras, including proofs of all the assertions, is beyond the scope of this volume. However, we shall include the definition of Borcherds algebras and the statements of the main results about their structure and representation theory, but without detailed proofs. In fact many of the results are quite similar to those we have already obtained about Kac-Moody algebras. However, the theory of Borcherds algebras includes examples which are quite different from Kac-Moody algebras. The best known such example is the Monster Lie algebra, which we shall describe in Section 21.3.

We begin with the definition of a Borcherds algebra. A Lie algebra $L$ over $\mathbb{R}$ is called a Borcherds algebra if it satisfies the following four axioms:
(i) $L$ has a $\mathbb{Z}$-grading

$$
L=\bigoplus_{i \in \mathbb{Z}} L_{i}
$$

such that $\operatorname{dim} L_{i}$ is finite for all $i \neq 0$, and $L$ is diagonalisable with respect to $L_{0}$. (Note that $\operatorname{dim} L_{0}$ need not be finite.)
(ii) There exists an automorphism $\omega: L \rightarrow L$ such that

$$
\begin{array}{ll}
\omega^{2}=1 & \\
\omega\left(L_{i}\right)=L_{-i} & \text { for all } i \in \mathbb{Z} \\
\omega=-1 & \text { on } L_{0} / L_{0} \cap Z(L)
\end{array}
$$

(iii) There is an invariant bilinear form

$$
\langle,\rangle: L \times L \rightarrow \mathbb{R}
$$

such that

$$
\begin{array}{ll}
\langle x, y\rangle=0 & \text { if } x \in L_{i}, y \in L_{j} \text { and } i+j \neq 0 \\
\langle w x, w y\rangle=\langle x, y\rangle & \text { for all } x, y \in L \\
-\langle x, w x\rangle>0 & \text { for } x \in L_{i} \text { with } i \neq 0, x \neq 0 .
\end{array}
$$

(iv) $L_{0} \subset[L L]$.

We observe some consequences of these axioms. In the first place it can be shown that $\left[L_{0} L_{0}\right]=0$, that is $L_{0}$ is abelian.

To describe a further consequence we define, for $x, y \in L$,

$$
\langle x, y\rangle_{0}=-\langle x, \omega y\rangle
$$

The scalar product $\langle,\rangle_{0}: L \times L \rightarrow \mathbb{R}$ is called the contravariant bilinear form. We now restrict the contravariant form to one of the graded components $L_{i}$ with $i \neq 0$ and have $\langle,\rangle_{0}: L_{i} \times L_{i} \rightarrow \mathbb{R}$. Let $x \in L_{i}$. Then

$$
\langle x, x\rangle_{0}=-\langle x, w x\rangle>0 \quad \text { if } x \neq 0
$$

Thus the contravariant form is positive definite on each graded component $L_{i}$ for $i \neq 0$ (though not necessarily on $L_{0}$ ).

We now give some examples of Borcherds algebras. We begin with a symmetric matrix $\mathfrak{H}$ over $\mathbb{R}$ which is either finite or countable. Thus

$$
\mathfrak{H}=\left(a_{i j}\right) \quad i, j \in I
$$

with $a_{i j} \in \mathbb{R}$ and $I$ either finite or countable. We shall assume that the matrix $\mathfrak{A}$ satisfies the conditions

$$
\begin{aligned}
& a_{i j} \leq 0 \quad \text { if } i \neq j \\
& \text { if } a_{i i}>0 \text { then } 2 a_{i j} / a_{i i} \in \mathbb{Z} \quad \text { for all } j .
\end{aligned}
$$

Proposition 21.1 There is a Borcherds algebra L associated to a symmetric matrix $\mathfrak{H}$ satisfying the above conditions which is defined as follows by generators and relations.

$$
L \text { is generated by elements } e_{i}, f_{i}, h_{i j} \quad i, j \in I
$$

subject to relations

$$
\begin{aligned}
& {\left[e_{i} f_{j}\right]=h_{i j}} \\
& {\left[h_{i j} h_{k l}\right]=0} \\
& {\left[h_{i j} e_{k}\right]=\delta_{i j} a_{i k} e_{k}} \\
& {\left[h_{i j} f_{k}\right]=-\delta_{i j} a_{i k} f_{k}}
\end{aligned}
$$

$$
\begin{aligned}
& \left(\operatorname{ad} e_{i}\right)^{n} e_{j}=0, \quad\left(\operatorname{ad} f_{i}\right)^{n} f_{j}=0 \quad \text { if } a_{i i}>0, \quad i \neq j \text { and } n=1-2 a_{i j} / a_{i i} \\
& {\left[e_{i} e_{j}\right]=0, \quad\left[f_{i} f_{j}\right]=0 \quad \text { if } a_{i i} \leq 0, \quad a_{j j} \leq 0 \text { and } a_{i j}=0}
\end{aligned}
$$

The Borcherds algebra $L$ defined by generators and relations in this way is called the universal Borcherds algebra associated with the symmetric matrix $\mathfrak{A}$. Its structure as a Borcherds algebra can be described as follows. Its involutary automorphism $\omega$ is given by

$$
\omega\left(e_{i}\right)=-f_{i}, \quad \omega\left(f_{i}\right)=-e_{i}, \quad \omega\left(h_{i j}\right)=-h_{j i}
$$

Its invariant bilinear form is uniquely determined by the condition

$$
\left\langle e_{i}, f_{i}\right\rangle=1 \quad \text { for all } i \in I
$$

In particular, if we write $h_{i}=h_{i i}$, then $\left[e_{i} f_{i}\right]=h_{i}$ and

$$
\left\langle h_{i}, h_{j}\right\rangle=\left\langle\left[e_{i} f_{i}\right], h_{j}\right\rangle=\left\langle e_{i},\left[f_{i} h_{j}\right]\right\rangle=\left\langle e_{i}, a_{i j} f_{i}\right\rangle=a_{i j}
$$

Thus

$$
\left\langle h_{i}, h_{j}\right\rangle=a_{i j} \text { for all } i, j \in I
$$

There are many ways of defining an appropriate grading on this Borcherds algebra. For each $i \in \mathbb{Z}$ let $n_{i} \in \mathbb{Z}$ satisfy $n_{i}>0$. Then there is a $\mathbb{Z}$-grading on $L$ uniquely determined by the conditions

$$
e_{i} \in L_{n_{i}}, \quad f_{i} \in L_{-n_{i}} .
$$

Further examples of Borcherds algebras can be obtained from a universal Borcherds algebra as follows. The axiom $\left[h_{i j} h_{k l}\right]=0$ shows that the subalgebra generated by all elements $h_{i j}, \quad i, j \in I$, is abelian. If $i \neq j$ then $h_{i j}$ lies in the centre of $L$ since $\left[h_{i j}, e_{k}\right]=0$ and $\left[h_{i j}, f_{k}\right]=0$ for all $k \in I$. Thus the subalgebra generated by the $h_{i j}$ for $i \neq j$ lies in the centre. It can be shown that the centre $Z$ of $L$ satisfies

$$
\left\langle h_{i j} ; \quad i, j \in I, i \neq j\right\rangle \subset Z \subset\left\langle h_{i j} ; \quad i, j \in I\right\rangle .
$$

In fact the Jacobi identity

$$
\left[\left[e_{i} f_{j}\right] h_{k}\right]+\left[\left[f_{j} h_{k}\right] e_{i}\right]+\left[\left[h_{k} e_{i}\right] f_{j}\right]=0
$$

shows that

$$
\left[h_{k},\left[e_{i} f_{j}\right]\right]=\left(a_{k i}-a_{k j}\right)\left[e_{i} f_{j}\right]
$$

It can be shown that, as a consequence of this, $h_{i j}=0$ unless $a_{k i}=a_{k j}$ for all $k \in I$, i.e. unless the $i$ th and $j$ th columns of $\mathfrak{H}$ are identical.

Proposition 21.2 Let L be a universal Borcherds algebra and I be an ideal of $L$ contained in the centre $Z$ of $L$. Then $L / I$ retains the structure of a Borcherds algebra.

The $\mathbb{Z}$-grading, involutary automorphism and invariant bilinear form on $L / I$ are readily obtained from those on $L$.

We now obtain still further Borcherds algebras. Starting from a universal Borcherds algebra $L$ we factor out an ideal $I$ of $L$ contained in the centre $Z$ of $L$. Then $L / I$ is still a Borcherds algebra. We write $\bar{L}=L / I$.

An inner derivation of $\bar{L}$ is one of form $x \rightarrow[x y]$ for some $y \in \bar{L}$, and an outer derivation is a derivation which either is zero or is not an inner derivation Let

$$
\bar{L}^{*}=\operatorname{Hom}(\bar{L}, \mathbb{R})
$$

and $A \subset \bar{L}^{*}$ be an abelian Lie algebra of outer derivations of $\bar{L}$. We suppose also that

$$
\left[\bar{e}_{i} x\right] \in \mathbb{R} \bar{e}_{i}, \quad\left[\bar{f}_{i} x\right] \in \mathbb{R} \bar{f}_{i}
$$

for all $x \in A$ where $\bar{e}_{i}, \bar{f}_{i}$ are images of $e_{i}, f_{i}$ under the natural homomorphism $L \rightarrow \bar{L}$.

Proposition 21.3 Let L be a universal Borcherds algebra and I be an ideal of $L$ contained in the centre $Z$ of $L$. Let $\bar{L}=L / I$. Let $A$ be an abelian Lie algebra of outer derivations of $\bar{L}$ and let $\bar{L}+A$ be the semidirect product of $\bar{L}$ by $A$ whose elements have form $x+a$ with $x \in \bar{L}, a \in A$ where

$$
[x+a, y+b]=[x y]+a(y)-b(x)
$$

Suppose that

$$
\left[\bar{e}_{i} x\right] \in \mathbb{R} \bar{e}_{i}, \quad\left[\bar{f}_{i} x\right] \in \mathbb{R} \bar{f}_{i}
$$

for all $x \in A$. Then $\bar{L}+A$ retains the structure of a Borcherds algebra in which $A \subset(\bar{L}+A)_{0}$.

The $\mathbb{Z}$-grading, involutary automorphism and invariant bilinear form on $\bar{L}+A$ are easily obtained from those of $\bar{L}$.

We have now constructed a family of Borcherds algebras which includes all universal Borcherds algebras, all quotients of such by ideals contained in the centre, and all semidirect products of such quotients by an abelian Lie algebra of outer derivations with suitable properties.

This turns out to give all possible Borcherds algebras, as is shown by the next theorem.

Theorem 21.4 Let L be a Borcherds algebra. Then there is a unique universal Borcherds algebra $L_{u}$ and a homomorphism

$$
f: L_{u} \rightarrow L
$$

(not necessarily unique) such that $\operatorname{ker} f$ is an ideal in the centre of $L_{u}, \operatorname{im} f$ is an ideal of $L$, and $L$ is the semidirect product of $\operatorname{im} f$ with an abelian Lie algebra of outer derivations lying in the 0 -graded component of $L$ and preserving all subspaces $\mathbb{R} \bar{e}_{i}$ and $\mathbb{R} \bar{f}_{i}$.

The homomorphism $f$ preserves the grading, involutary automorphism, and bilinear form.

Now that we have obtained the complete set of Borcherds algebras in this way, we explore their relationship with symmetrisable Kac-Moody algebras. It turns out that every symmetrisable Kac-Moody algebra over $\mathbb{R}$ gives rise to a universal Borcherds algebra, which is the subalgebra of the Kac-Moody algebra obtained by generators and relations prior to the extension of the Cartan subalgebra by an abelian Lie algebra of outer derivations.

Theorem 21.5 Let L be a symmetrisable Kac-Moody algebra with GCM $A=\left(A_{i j}\right)$. Thus there exists a diagonal matrix $D=\operatorname{diag}\left(d_{1}, \ldots, d_{n}\right)$ with each $d_{i} \in \mathbb{Z}, d_{i}>0$ such that $D A$ is symmetric. Let $\mathfrak{A}=\left(a_{i j}\right)$ be given by

$$
a_{i j}=\frac{d_{i} A_{i j}}{2}
$$

Then we have $a_{i j}=a_{j i}$ and $a_{i i}=d_{i}$. Thus $a_{i j} \leq 0$ if $i \neq j$ and $a_{i i}$ is a positive integer. Moreover $2 a_{i j} / a_{i i}=A_{i j} \in \mathbb{Z}$.

Then the symmetric matrix $\left(a_{i j}\right)$ satisfies the conditions needed to construct a Borcherds algebra, and the universal Borcherds algebra with symmetric matrix $\mathfrak{H}$ coincides with the subalgebra of the Kac-Moody algebra L obtained by generators and relations prior to the adjunction of the abelian Lie algebra of outer derivations.

In this way every symmetrisable Kac-Moody algebra determines a certain subalgebra which is a universal Borcherds algebra. In fact the main points of difference between symmetrisable Kac-Moody algebras and universal Borcherds algebras are that, in a Borcherds algebra:

I may be countably infinite rather than finite
$a_{i i}$ may not be positive and need not lie in $\mathbb{Z}$
$2 a_{i j} / a_{i i}$ is only assumed to lie in $\mathbb{Z}$ when $a_{i i}>0$.

### 21.2 Representations of Borcherds algebras

We now introduce the root system and Weyl group of a Borcherds algebra.
We suppose first that $L$ is a universal Borcherds algebra. The root lattice $Q$ of $L$ is the free abelian group with basis $\alpha_{i}$ for $i \in I$. We have a symmetric bilinear form

$$
Q \times Q \rightarrow \mathbb{R}
$$

given by $\left(\alpha_{i}, \alpha_{j}\right) \rightarrow\left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i j}$.
The basis elements $\alpha_{i}$ of $Q$ are called the fundamental roots. The set of fundamental roots is denoted by $\Pi$. We have a grading

$$
L=\bigoplus_{\alpha \in Q} L_{\alpha}
$$

determined by the conditions

$$
e_{i} \in L_{\alpha_{i}}, \quad f_{i} \in L_{-\alpha_{i}}
$$

An element $\alpha \in Q$ is called a root of $L$ if $\alpha \neq 0$ and $L_{\alpha} \neq O . \alpha$ is called a positive root if $\alpha$ is a sum of fundamental roots. For any root $\alpha$ either $\alpha$ or $-\alpha$ is positive. We have

$$
\Phi=\Phi^{+} \cup \Phi^{-}
$$

where $\Phi$ is the set of roots and $\Phi^{+}, \Phi^{-}$are the subsets of positive and negative roots respectively. We say that $\alpha \in \Phi$ is a real root if $\langle\alpha, \alpha\rangle>0$ and $\alpha \in \Phi$ is an imaginary root if $\langle\alpha, \alpha\rangle \leq 0$.

We next introduce the Weyl group $W$ of the universal Borcherds algebra $L$. $W$ is the group of isometries of the root lattice $Q$ generated by the reflections $s_{i}$ corresponding to the real fundamental roots. We have

$$
s_{i}\left(\alpha_{j}\right)=\alpha_{j}-2 \frac{\left\langle\alpha_{i}, \alpha_{j}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} \alpha_{i}=\alpha_{j}-2 \frac{a_{i j}}{a_{i i}} \alpha_{i}
$$

We recall that $2 a_{i j} / a_{i i} \in \mathbb{Z}$ since $a_{i i}>0$.

Let $H$ be the abelian subalgebra of $L$ generated by the elements $h_{i j}$ for all $i, j \in I$. We have a map

$$
Q \rightarrow H
$$

under which $\alpha_{i}$ maps to $h_{i}$, which is a homomorphism of abelian groups and preserves the scalar product. However, this map need not necessarily be injective.

So for we have assumed that $L$ is a universal Borcherds algebra. However, if $L$ is an arbitrary Borcherds algebra there is an associated universal Borcherds algebra $L_{u}$ given by Theorem 21.4. Then the root system of $L$ is defined to be the root system of $L_{u}$ and the Weyl group of $L$ is defined to be the Weyl group of $L_{u}$.

This theory of Borcherds algebras is thus very similar to the theory of Kac-Moody algebras. The main difference is that for Borcherds algebras there can exist imaginary fundamental roots, and that the Weyl group is generated by the reflections with respect to the real fundamental roots only.

We now turn to the representation theory of Borcherds algebras. We define the set $X$ of integral weights by

$$
X=\left\{\lambda \in Q \otimes \mathbb{R} ; 2 \frac{\left\langle\lambda, \alpha_{i}\right\rangle}{\left\langle\alpha_{i}, \alpha_{i}\right\rangle} \in \mathbb{Z} \quad \text { for all } \alpha_{i} \in \Pi_{\mathrm{Re}}\right\}
$$

Here $\Pi_{\mathrm{Re}}$ is the set of real fundamental roots. We recall that $\left\langle\alpha_{i}, \alpha_{i}\right\rangle>0$ when $\alpha_{i} \in \Pi_{\mathrm{Re}}$. We define the subset $X^{+} \subset X$ of dominant integral weights by

$$
X^{+}=\left\{\lambda \in X ;\left\langle\lambda, \alpha_{i}\right\rangle \geq 0 \quad \text { for all } \alpha_{i} \in \Pi\right\}
$$

In a manner very similar to that we have described for Kac-Moody algebras in Chapter 19 it is possible to define an irreducible module $L(\lambda)$ for the Borcherds algebra $L$ associated to any dominant integral weight $\lambda . L(\lambda)$ is called the irreducible $L$-module with highest weight $\lambda$.

Now Borcherds proved a character formula for $L(\lambda)$ analogous to Kac' character formula Theorem 19.16 for Kac-Moody algebras.

Theorem 21.6 (Borcherds' character formula). Let L be a Borcherds algebra, $\lambda$ a dominant integral weight and $L(\lambda)$ the corresponding irreducible $L$-module with highest weight $\lambda$. Then the character of $L(\lambda)$ is given by

$$
\operatorname{ch} L(\lambda)=\frac{\sum_{w \in W} \varepsilon(w) w\left(\sum_{\Psi}(-1)^{|\Psi|} e\left(\lambda+\rho-\sum \Psi\right)\right)}{e(\rho) \prod_{\alpha \in \Phi^{+}}(1-e(-\alpha))^{m_{\alpha}}}
$$

where $m_{\alpha}=\operatorname{dim} L_{\alpha}, \Psi$ runs over all finite subsets of mutually orthogonal imaginary fundamental roots, and $\rho$ is any element of $Q \otimes \mathbb{R}$ such that

$$
\left\langle\rho, \alpha_{i}\right\rangle=\frac{1}{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle
$$

for all real fundamental roots $\alpha_{i}$.

As usual this character is interpreted as

$$
\operatorname{ch} L(\lambda)=\sum_{\mu}\left(\operatorname{dim} L(\lambda)_{\mu}\right) e(\mu)
$$

where $\mu \rightarrow e(\mu)$ is an isomorphism between the additive group of weights and the corresponding multiplicative group.
(In fact there may not exist a vector $\rho \in Q \otimes \mathbb{R}$ such that $\left\langle\rho, \alpha_{i}\right\rangle=\frac{1}{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle$ for all real fundamental roots $\alpha_{i}$ of a general Borcherds algebra. But if there is no such $\rho \in Q \otimes \mathbb{R}, \rho$ may still be defined as the homomorphism from $Q$ to $\mathbb{R}$ taking $\alpha_{i}$ to $\frac{1}{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle$ for all $i \in I$, and the character formula can be interpreted accordingly.)

In the special case $\lambda=0, L(\lambda)$ is the trivial 1-dimensional $L$-module. Then Borcherds' character formula reduces to the following identity.

Theorem 21.7 (Borcherds' denominator formula).

$$
e(\rho) \prod_{\alpha \in \Phi^{+}}(1-e(-\alpha))^{m_{\alpha}}=\sum_{w \in W} \varepsilon(w) w\left(e(\rho) \sum_{\Psi}(-1)^{|\Psi|} e\left(-\sum \Psi\right)\right)
$$

By substituting Borcherds' denominator formula into Theorem 21.6 we obtain a second form of Borcherds' character formula.

Theorem 21.8 (Borcherds' character formula, second form). With the notation of Theorem 21.6 we have

$$
\operatorname{ch} L(\lambda)=\frac{\sum_{w \in W} \varepsilon(w) w\left(\sum_{\Psi}(-1)^{|\Psi|} e\left(\lambda+\rho-\sum \Psi\right)\right)}{\sum_{w \in W} \varepsilon(w) w\left(e(\rho) \sum_{\Psi}(-1)^{|\Psi|} e\left(-\sum \Psi\right)\right)}
$$

## Comments on the proof of Borcherds' character formula

We shall not give the proof of Borcherds' character formula in detail, since the ideas are quite similar to those which arise in the proof of Kac' character
formula for Kac-Moody algebras. However we shall say enough to explain where the additional term

$$
\sum_{\Psi}(-1)^{|\Psi|} e\left(-\sum \Psi\right)
$$

comes from, where $\Psi$ runs over all finite sets of mutually orthogonal imaginary fundamental roots. Of course in Kac-Moody algebras there are no imaginary fundamental roots so the only possible subset $\Psi$ is the empty set. The additional term then becomes $e(0)$, the identity element of $e(Q)$, and disappears from the formula.

For a Borcherds algebra $L$ we have

$$
\Pi=\Pi_{\mathrm{Re}} \cup \Pi_{\mathrm{Im}}
$$

where $\Pi_{\mathrm{Re}}$ is the set of real fundamental roots and $\Pi_{\mathrm{Im}}$ is the set of imaginary fundamental roots. We also define

$$
\Phi_{\mathrm{Re}}=W\left(\Pi_{\mathrm{Re}}\right), \quad \Phi_{\mathrm{Im}}=\Phi-\Phi_{\mathrm{Re}}
$$

to be the sets of real and imaginary roots respectively. We recall from Theorem 16.24 that, in a Kac-Moody algebra,

$$
\Phi_{\mathrm{Im}}^{+}=\bigcup_{w \in W} w(K)
$$

where $K=\left\{\alpha \in Q^{+} ; \alpha \neq 0, \operatorname{supp} \alpha\right.$ is connected, $\left.-\alpha \in \bar{C}\right\}$ and

$$
\bar{C}=\left\{\lambda \in Q \otimes \mathbb{R} ;\left\langle\lambda, \alpha_{i}\right\rangle \geq 0 \text { for all } \alpha_{i} \in \Pi_{\mathrm{Re}}\right\} .
$$

There is an analogous result for Borcherds algebras given as follows.

Theorem 21.9 The set of positive imaginary roots of a Borcherds algebra is given by

$$
\Phi_{\mathrm{Im}}^{+}=\bigcup_{w \in W} w(K)
$$

where $K$ is given by

$$
\begin{aligned}
K= & \left\{\alpha \in Q^{+} ; \alpha \neq 0, \quad-\alpha \in \bar{C}, \quad \operatorname{supp} \alpha \text { is connected }\right\} \\
& -\left\{j \alpha_{i} ; j \in \mathbb{Z}, \quad j \geq 2, \quad \alpha_{i} \in \Pi_{\mathrm{Im}}\right\}
\end{aligned}
$$

Proof. Omitted. The idea is generally similar to that of Theorem 16.24. It is clearly necessary to exclude positive multiples $j \alpha_{i}$ of imaginary fundamental
roots with $j \geq 2$ since these vectors satisfy the conditions required for belonging to $K$, but cannot be roots since there is no possible root vector giving rise to such a root.

Following closely the proof of Kac' character formula we obtain

$$
e(\rho) \prod_{\alpha \in \Phi^{+}}(1-e(-\alpha))^{m_{\alpha}} \operatorname{ch} L(\lambda)=\sum_{\mu} c_{\mu} e(\mu+\rho)
$$

summed over all weights $\mu$ such that $\mu \prec \lambda$ and $\langle\mu+\rho, \mu+\rho\rangle=$ $\langle\lambda+\rho, \lambda+\rho\rangle$, where $c_{\mu} \in \mathbb{Z}$ and both sides are skew-symmetric under the action of the Weyl group $W$. (See the proof of Theorem 19.16.)

We now define a certain partial sum $S$ of terms on the right-hand side.
Let $S=\sum_{\mu} c_{\mu} e(\mu+\rho)$, summed over all weights $\mu$ satisfying $\mu \prec \lambda$, $\langle\mu+\rho, \mu+\rho\rangle=\langle\lambda+\rho, \lambda+\rho\rangle$ and $\left\langle\mu+\rho, \alpha_{i}\right\rangle \geq 0$ for all $\alpha_{i} \in \Pi_{\mathrm{Re}}$.

Since $\mu \prec \lambda$ we have

$$
\mu=\lambda-\sum k_{i} \alpha_{i} \quad \text { for some } k_{i} \in \mathbb{Z}, k_{i}>0, \alpha_{i} \in \Pi
$$

Since $\langle\mu+\rho, \mu+\rho\rangle=\langle\lambda+\rho, \lambda+\rho\rangle$ we have

$$
\langle(\lambda+\rho)-(\mu+\rho),(\lambda+\rho)+(\mu+\rho)\rangle=0
$$

that is $\left\langle\sum k_{i} \alpha_{i}, \lambda+\mu+2 \rho\right\rangle=0$. This implies

$$
\sum k_{i}\left\langle\alpha_{i}, \lambda\right\rangle+\sum k_{i}\left\langle\alpha_{i}, \mu+2 \rho\right\rangle=0
$$

We can deduce several consequences from this equation. We note first that $\left\langle\alpha_{i}, \lambda\right\rangle \geq 0$ since $\lambda$ is dominant. Also for $\alpha_{i} \in \Pi_{\mathrm{Re}}$ we have

$$
\left\langle\alpha_{i}, \mu+2 \rho\right\rangle=\left\langle\alpha_{i}, \mu+\rho\right\rangle+\left\langle\alpha_{i}, \rho\right\rangle=\left\langle\alpha_{i}, \mu+\rho\right\rangle+\frac{1}{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle>\left\langle\alpha_{i}, \mu+\rho\right\rangle .
$$

Now $\left\langle\alpha_{i}, \mu+\rho\right\rangle \geq 0$ by definition of $S$, so $\left\langle\alpha_{i}, \mu+2 \rho\right\rangle>0$. On the other hand, for $\alpha_{i} \in \Pi_{\mathrm{Im}}$ we have

$$
\left\langle\alpha_{i}, \mu+2 \rho\right\rangle=\left\langle\alpha_{i}, \mu+\alpha_{i}\right\rangle=\left\langle\alpha_{i}, \lambda-\sum k_{j}^{\prime} \alpha_{j}\right\rangle \quad \text { for some } k_{j}^{\prime}>0
$$

Thus $\left\langle\alpha_{i}, \mu+2 \rho\right\rangle \geq 0$ since $\left\langle\alpha_{i}, \lambda\right\rangle \geq 0,\left\langle\alpha_{i}, \alpha_{j}\right\rangle \leq 0$ if $i \neq j$, and $\left\langle\alpha_{i}, \alpha_{i}\right\rangle \leq 0$. We now collect these results together and put them into the equation

$$
\sum k_{i}\left\langle\alpha_{i}, \lambda\right\rangle+\sum k_{i}\left\langle\alpha_{i}, \mu+2 \rho\right\rangle=0
$$

The conclusion is that $\left\langle\alpha_{i}, \lambda\right\rangle=0$ and $\left\langle\alpha_{i}, \mu+2 \rho\right\rangle=0$ for all $\alpha_{i}$ in the sum $\sum k_{i} \alpha_{i}$. This in turn implies that each such $\alpha_{i} \in \Pi_{\mathrm{Im}}$. But then

$$
\left\langle\alpha_{i}, \mu+2 \rho\right\rangle=\left\langle\alpha_{i}, \lambda\right\rangle-\sum k_{j}^{\prime}\left\langle\alpha_{i}, \alpha_{j}\right\rangle=-\sum k_{j}^{\prime}\left\langle\alpha_{i}, \alpha_{j}\right\rangle
$$

Since $k_{j}^{\prime}>0$ and $\left\langle\alpha_{i}, \alpha_{j}\right\rangle \leq 0$ this implies that $\left\langle\alpha_{i}, \alpha_{j}\right\rangle=0$. In fact $k_{j}^{\prime}=k_{j}$ if $j \neq i$ and $k_{i}^{\prime}=k_{i}-1$. So if $i \neq j$ we have $\left\langle\alpha_{i}, \alpha_{j}\right\rangle=0$ for all $\alpha_{i}, \alpha_{j}$ in the sum $\sum k_{i} \alpha_{i}$ with $i \neq j$. In other words,

$$
\lambda-\mu=\sum k_{i} \alpha_{i}
$$

is a linear combination of mutually orthogonal imaginary fundamental roots all of which are orthogonal to $\lambda$.

Now we have

$$
\left\langle\mu, \alpha_{j}\right\rangle=\left\langle\lambda, \alpha_{j}\right\rangle-\sum_{i} k_{i}\left\langle\alpha_{i}, \alpha_{j}\right\rangle \geq 0
$$

for all $\alpha_{j} \in \Pi_{\mathrm{Re}}$ since $\left(\lambda, \alpha_{j}\right) \geq 0, k_{i}>0$ and $\left\langle\alpha_{i}, \alpha_{j}\right\rangle \leq 0$ since $\alpha_{i} \in \Pi_{\mathrm{Im}}, \alpha_{j} \in \Pi_{\mathrm{Re}}$ so $i \neq j$. Thus $\mu \in \bar{C}$. It follows that $\mu+\rho \in C$ since

$$
\left\langle\mu+\rho, \alpha_{j}\right\rangle=\left\langle\mu, \alpha_{j}\right\rangle+\frac{1}{2}\left\langle\alpha_{j}, \alpha_{j}\right\rangle>\left\langle\mu, \alpha_{j}\right\rangle
$$

so $\left\langle\mu+\rho, \alpha_{j}\right\rangle>0$. Thus $\mu+\rho$ lies in the fundamental chamber $C$. Since our sum

$$
\sum_{\substack{\mu \\ \mu \lambda \\\langle\mu+\rho, \mu+\rho\rangle=\langle\lambda+\rho, \lambda+\rho\rangle}} c_{\mu} e(\mu+\rho)
$$

is skew-symmetric under the action of $W$, this sum must be equal to

$$
\sum_{w \in W} \varepsilon(w) w(S)
$$

since $S$ is the partial sum including all summands $c_{\mu} e(\mu+\rho)$ for which $\mu+\rho$ lies in $\bar{C}$.

We shall now determine $S$. Let the module $L(\lambda)$ have highest weight vector $v_{\lambda}$. If $\lambda\left(h_{i}\right)=0$ then $f_{i} v_{\lambda}=0$ by the analogue in Borcherds algebras of the proof of Theorem 10.20. If all $\alpha_{i}$ in the sum $\sum k_{i} \alpha_{i}$ satisfied $\lambda\left(h_{i}\right)=0$ then we would have $f_{i} v_{\lambda}=0$ for all such $i$ and this would imply that $\lambda-\sum k_{i} \alpha_{i}$ could not be a weight of $L(\lambda)$. So if $\lambda-\sum k_{i} \alpha_{i}$ is a weight not equal to $\lambda$ we must have $\lambda\left(h_{i}\right) \neq 0$ for some $i$ in this sum. Thus $\left\langle\lambda, \alpha_{i}\right\rangle \neq 0$ for some $i$ in this sum. On the other hand we know that $\mu=\lambda-\sum k_{i} \alpha_{i}$ where all $\alpha_{i}$ in the sum satisfy $\left\langle\lambda, \alpha_{i}\right\rangle=0$. This implies that the weight

$$
\mu+\rho=(\lambda+\rho)-\sum k_{i} \alpha_{i}
$$

can only arise from the term $e(\lambda)$ in $\operatorname{ch} L(\lambda)$ in the formula

$$
e(\rho) \prod_{\alpha \in \Phi^{+}}(1-e(-\alpha))^{m_{\alpha}} \operatorname{ch} L(\lambda)=\sum_{\mu} e_{\mu} e(\mu+\rho) .
$$

So all terms on the right of this formula which lie in $S$ arise from

$$
e(\lambda+\rho) \prod_{\alpha \in \Phi^{+}}(1-e(-\alpha))^{m_{\alpha}}
$$

on the left.
We consider which roots $\alpha \in \Phi^{+}$in the formula

$$
e(\lambda+\rho) \prod_{\alpha \in \Phi^{+}}(1-e(-\alpha))^{m_{\alpha}}
$$

can contribute to give $\mu+\rho=(\lambda+\rho)-\sum k_{i} \alpha_{i}$. All such roots $\alpha \in \Phi^{+}$must be linear combinations of the fundamental roots $\alpha_{i}$ arising in the sum $\sum k_{i} \alpha_{i}$. But such $\alpha_{i}$ are mutually orthogonal imaginary fundamental roots. So sums of two or more such $\alpha_{i}$ do not have connected support, so cannot be roots by Theorem 21.9. Also each $\alpha_{i} \in \Pi$ has $m_{\alpha_{i}}=1$ since the corresponding root space is spanned by $e_{i}$. Thus a weight $\mu+\rho$ giving a term on the right which lies in $S$ must arise from

$$
e(\lambda+\rho) \prod_{\alpha_{i} \in \Pi_{\mathrm{Im}}}\left(1-e\left(-\alpha_{i}\right)\right)=e(\lambda+\rho) \sum_{\Psi}(-1)^{|\Psi|} e\left(-\sum \Psi\right)
$$

summed over all finite sets $\Psi$ of mutually orthogonal imaginary fundamental roots. Thus we have

$$
S=\sum_{\Psi}(-1)^{|\Psi|} e\left(\lambda+\rho-\sum \Psi\right)
$$

and so

$$
e(\rho) \prod_{\alpha \in \Phi^{+}}(1-e(-\alpha))^{m_{\alpha}} \operatorname{ch} L(\lambda)=\sum_{w \in W} \varepsilon(w) w\left(\sum_{\Psi}(-1)^{|\Psi|} e\left(\lambda+\rho-\sum \Psi\right)\right)
$$

as required.
This argument therefore explains the difference between Kac' character formula and Borcherds' character formula, and where the extra term in Borcherds' character formula comes from.

### 21.3 The Monster Lie algebra

In this final section we shall show that, although Borcherds algebras have many properties which seem quite similar to those of Kac-Moody algebras, they include examples which behave in a very different way from Kac-Moody algebras. The example we have in mind is the Monster Lie algebra. The definition and properties of the Monster Lie algebra are closely related to the properties of a certain modular function $j$, so we shall begin by describing the definition and significance of this function.

We first recall the action of the group $S L_{2}(\mathbb{R})$ on the upper half plane $H$. Let

$$
\begin{aligned}
& S L_{2}(\mathbb{R})=\left\{\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) ; a d-b c=1, \quad a, b, c, d \in \mathbb{R}\right\} \\
& H=\{\tau \in \mathbb{C} ; \operatorname{Im} \tau>0\}
\end{aligned}
$$

The group $S L_{2}(\mathbb{R})$ acts on $H$ by

$$
\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \tau=\frac{a \tau+b}{c \tau+d}
$$

since if $\operatorname{Im} \tau>0$ we have $\operatorname{Im}\left(\frac{a \tau+b}{c \tau+d}\right)>0$. In particular the subgroup $S L_{2}(\mathbb{Z})$ acts on $H$. Since

$$
\left(\begin{array}{cc}
-1 & 0 \\
0 & -1
\end{array}\right) \tau=\tau
$$

we see that $P S L_{2}(\mathbb{Z})=S L_{2}(\mathbb{Z}) /\left( \pm I_{2}\right)$ acts on $H . P S L_{2}(\mathbb{Z})$ is called the modular group.

We denote by $H / S L_{2}(\mathbb{Z})$ the set of orbits. Since $\left(\begin{array}{ll}1 & 1 \\ 0 & 1\end{array}\right) \in S L_{2}(\mathbb{Z})$ the elements $\tau$ and $\tau+1$ of $H$ lie in the same orbit. Thus each orbit intersects

$$
\left\{\tau \in H ;-\frac{1}{2} \leq \operatorname{Re} \tau \leq \frac{1}{2}\right\}
$$

Again we have $\left(\begin{array}{cc}0 & 1 \\ -1 & 0\end{array}\right) \in S L_{2}(\mathbb{Z})$ and so the elements $\tau,-1 / \tau \in H$ lie in the same orbit. Thus each orbit intersects

$$
\{\tau \in H ;|\tau| \geq 1\}
$$

In fact we can obtain a fundamental region for the action of $S L_{2}(\mathbb{Z})$ on $H$ by taking the region

$$
\left\{\tau \in H ;-\frac{1}{2} \leq \operatorname{Re} \tau \leq \frac{1}{2}, \quad|\tau| \geq 1\right\}
$$

and identifying the points $\tau, \tau+1$ for $\operatorname{Re} \tau=-\frac{1}{2}$ and the points $\sigma,-1 / \sigma$ for $|\sigma|=1$. The fundamental region is illustrated in Figure 21.1.

Having made the above identifications we obtain a set intersecting each orbit in just one point. The set $H / S L_{2}(\mathbb{Z})$ of orbits has the structure of a compact Riemann surface with one point removed. This is a Riemann surface of genus 0 , i.e. a Riemann sphere. When we remove one point from it we


Figure 21.1 Fundamental region
obtain a subset which can be identified with $\mathbb{C}$. Thus we have an isomorphism of Riemann surfaces

$$
H / S L_{2}(\mathbb{Z}) \rightarrow \mathbb{C}
$$

This can be extended to an isomorphism of compact Riemann surfaces by adding the point $\mathrm{i} \infty$ on the left and $\infty$ on the right. Thus we have an isomorphism

$$
\left(H / S L_{2}(\mathbb{Z})\right) \cup\{\mathrm{i} \infty\} \rightarrow S^{2}=\mathbb{C} \cup\{\infty\}
$$

under which im maps to $\infty$. Such an isomorphism of Riemann surfaces is not uniquely determined. However, if $j$ is any such isomorphism any other must have the form $a(j+b)$ where $a, b$ are constants and $a \neq 0$. Such a map determines a map from $H$ to $\mathbb{C}$ constant on orbits. This map will also be denoted by $j . j$ is a modular function, i.e. a function invariant under the action of the modular group.

Since $\tau, \tau+1$ lie in the same orbit we have $j(\tau)=j(\tau+1)$, thus $j$ is periodic. This implies that $j$ has a Fourier expansion of form

$$
j(\tau)=\sum_{n \in \mathbb{Z}} c_{n} \mathrm{e}^{2 \pi \mathrm{i} n \tau}
$$

We write $q=\mathrm{e}^{2 \pi i \tau}$. Then we have

$$
j(\tau)=\sum_{n \in \mathbb{Z}} c_{n} q^{n} .
$$

We shall now describe such a function $j$. In order to do so we first introduce some modular forms. A function $f: H \rightarrow \mathbb{C}$ is called a modular form of weight $k$ if

$$
f\left(\frac{a \tau+b}{c \tau+d}\right)=(c \tau+d)^{k} f(\tau)
$$

for all $\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in S L_{2}(\mathbb{Z})$. We give two examples of modular forms. The first is an example of a so-called Eisenstein series. For each positive integer $n$ let

$$
\sigma_{3}(n)=\sum_{d \mid n} d^{3}
$$

summed over all divisors of $n$, and let

$$
E_{4}(\tau)=1+240 \sum_{n \geq 1} \sigma_{3}(n) q^{n}
$$

Thus

$$
E_{4}(\tau)=1+240 q+2160 q^{2}+\cdots
$$

This function is known to be a modular form of weight 4.
Secondly define $\Delta$ by

$$
\Delta(\tau)=q \prod_{n \geq 1}\left(1-q^{n}\right)^{24}
$$

Then

$$
\Delta(\tau)=q-24 q^{2}+252 q^{3}-\cdots
$$

This is called Dedekind's $\Delta$-function and is known to be a modular form of weight 12 .

We now define $j: H \rightarrow \mathbb{C}$ by

$$
j(\tau)=\frac{E_{4}(\tau)^{3}}{\Delta(\tau)}
$$

This is a modular form of weight 0 , i.e. a modular function, and so is constant on orbits of $S L_{2}(\mathbb{Z})$ on $H$. We have

$$
j(\tau)=q^{-1}+744+196884 q+21493760 q^{2}+\cdots
$$

and $j$ has a simple pole at $\tau=\mathrm{i} \infty$, i.e. $q=0$.
$j$ gives an isomorphism of Riemann surfaces

$$
j: H / S L_{2}(\mathbb{Z}) \rightarrow \mathbb{C}
$$

which extends to an isomorphism of compact Riemann surfaces

$$
j:\left(H / S L_{2}(\mathbb{Z})\right) \cup\{\mathrm{i} \infty\} \rightarrow S^{2}=\mathbb{C} \cup\{\infty\}
$$

Any other such isomorphism has the form $a(j+b)$ where $a, b$ are constants and $a \neq 0$. In particular there is just one such isomorphism with leading
coefficient 1 and constant term 0 . We shall call this the canonical isomorphism. This is the function

$$
j(\tau)-744=q^{-1}+\sum_{n \geq 1} c_{n} q^{n}
$$

where $c_{1}=196884, c_{2}=21493760$, etc. All the $c_{n}$ are positive integers.
We are now ready to introduce the Monster Lie algebra. We first define a countable symmetric matrix $\mathfrak{A}$. $\mathfrak{H}$ is defined as a block matrix, with blocks of rows and columns parametrised by the natural numbers $\mathbb{N}=\{0,1,2,3, \ldots\}$. Let $B_{i j}$ be the $(i, j)$-block of $\mathfrak{A}$. The number of rows in $B_{i j}$ is 1 if $i=0$ and $c_{i}$ if $i \neq 0$ where $c_{i}$ is the coefficient of $q^{i}$ in the modular function $j$. Similarly the number of columns in $B_{i j}$ is 1 if $j=0$ and $c_{j}$ if $j \neq 0$. All the matrix entries in a given block $B_{i j}$ are equal to one another. These entries are given as follows.

The single entry in block $B_{00}$ is 2 .
All entries in block $B_{0 n}$ for $n \neq 0$ are $-(n-1)$.
All entries in block $B_{m n}$ for $m \neq 0, n \neq 0$ are $-(m+n)$.
These conditions determine the matrix $\mathfrak{H}$. We have

Let $L(\mathfrak{H})$ be the universal Borcherds algebra determined by the countable matrix $\mathfrak{A}$ as in Proposition 21.1. Let $Z$ be the centre of $L(\mathfrak{H})$ and $H$ be the subalgebra of $L(\mathfrak{H})$ generated by all elements $h_{i j}=\left[e_{i} f_{j}\right]$. We know from Section 21.1 that

$$
\left\langle h_{i j} ; i \neq j\right\rangle \subset Z \subset H
$$

It is also clear that $h_{i}-h_{j} \in Z$ where $h_{i}=h_{i i}, h_{j}=h_{j j}$ and $i, j$ are in the same block, since columns $i, j$ of $\mathfrak{H}$ are then identical. $Z$ is in fact generated by the
elements $h_{i j}$ for $i \neq j$ and $h_{i}-h_{j}$ where $i, j$ are in the same block. Moreover $Z$ is an ideal of $L(\mathfrak{H})$.

Let $\mathfrak{M}=L(\mathfrak{H}) / Z . \mathfrak{M}$ is called the Monster Lie algebra. We shall now determine some properties of $\mathfrak{M}$.

Let the blocks of rows of $\mathfrak{H}$ be $B_{0}, B_{1}, B_{2}, \ldots$ with $\left|B_{0}\right|=1,\left|B_{n}\right|=c_{n}$ for $n \geq 1$. We choose one $i \in I$ out of each block, such that $i$ lies in the block $B_{i}$. We then consider the elements $h_{i}$ for such elements $i \in I$. Thus we have elements $h_{0}, h_{1}, h_{2}, \ldots \in H$. Then the elements

$$
\begin{aligned}
& 2 h_{2}+h_{0}-3 h_{1} \\
& 2 h_{3}+2 h_{0}-4 h_{1} \\
& 2 h_{4}+3 h_{0}-5 h_{1}
\end{aligned}
$$

all lie in $Z$, since the corresponding linear combinations of the rows of $\mathfrak{A}$ are all zero vectors.

Let $h_{i} \rightarrow \bar{h}_{i}$ under the natural homomorphism $L(\mathfrak{H}) \rightarrow \mathfrak{M}$. Then we have

$$
\begin{aligned}
& 2 \bar{h}_{2}=3 \bar{h}_{1}-\bar{h}_{0} \\
& 2 \bar{h}_{3}=4 \bar{h}_{1}-2 \bar{h}_{0} \\
& 2 \bar{h}_{4}=5 \bar{h}_{1}-3 \bar{h}_{0}
\end{aligned}
$$

Let $\mathfrak{M}_{0}=H / Z . \mathfrak{M}_{0}$ is the Cartan subalgebra of $\mathfrak{M}$, being the image of $H$ under the natural homomorphism. We see from the above relations that $\mathfrak{M}_{0}$ is spanned by $\bar{h}_{0}$ and $\bar{h}_{1}$. Moreover $\bar{h}_{0}$ and $\bar{h}_{1}$ are linearly independent since this is true of the first two rows of $\mathfrak{M}$. Thus $\bar{h}_{0}, \bar{h}_{1}$ form a basis of $\mathfrak{M}_{0}$ and we have

$$
\operatorname{dim} \mathfrak{M}_{0}=2
$$

In fact we find it more convenient to choose the basis $b_{0}, b_{1}$ of $\mathfrak{M}_{0}$ given by

$$
b_{0}=\frac{\bar{h}_{0}+\bar{h}_{1}}{2}, \quad b_{1}=\frac{-\bar{h}_{0}+\bar{h}_{1}}{2}
$$

Thus $\mathfrak{M}_{0}=\mathbb{R} b_{0}+\mathbb{R} b_{1}$. The scalar product on $\mathfrak{M}_{0}$ is given by

$$
\begin{aligned}
& \left\langle\bar{h}_{0}, \bar{h}_{0}\right\rangle=a_{00}=2 \\
& \left\langle\bar{h}_{1}, \bar{h}_{1}\right\rangle=a_{11}=-2 \\
& \left\langle\bar{h}_{0}, \bar{h}_{1}\right\rangle=\left\langle\bar{h}_{1}, \bar{h}_{0}\right\rangle=a_{01}=0 .
\end{aligned}
$$

It follows that

$$
\left\langle b_{0}, b_{0}\right\rangle=0, \quad\left\langle b_{1}, b_{1}\right\rangle=0, \quad\left\langle b_{0}, b_{1}\right\rangle=-1
$$

Hence

$$
\left\langle m b_{0}+n b_{1}, m^{\prime} b_{0}+n^{\prime} b_{1}\right\rangle=-\left(m n^{\prime}+n m^{\prime}\right) .
$$

We now regard the Monster Lie algebra $\mathfrak{M}$ as a module over its Cartan subalgebra $\mathfrak{M}_{0}$. Let $m, n \in \mathbb{Z}$ and define $\mathfrak{M}_{(m, n)}$ by

$$
\mathfrak{M}_{(m, n)}=\left\{x \in \mathfrak{M} ;\left[b_{0} x\right]=m b_{0}, \quad\left[b_{1} x\right]=n b_{1}\right\}
$$

Then one can show that $\mathfrak{M}_{(0,0)}=\mathfrak{M}_{0}$ and

$$
\mathfrak{M}=\bigoplus \mathcal{M}_{(m, n)} \quad \text { for }(m, n) \in \mathbb{Z} \times \mathbb{Z}
$$

Moreover we have

$$
\begin{aligned}
& \operatorname{dim} \mathfrak{M}_{(m, n)}=c_{m n} \quad \text { if } m \neq 0, n \neq 0 \\
& \operatorname{dim} \mathfrak{M}_{(0,0)}=2 \\
& \operatorname{dim} \mathfrak{M}_{(m, 0)}=\operatorname{dim} \mathfrak{M}_{(0, n)}=0 \quad \text { if } m \neq 0, n \neq 0
\end{aligned}
$$

(These results follow from the 'no-ghost' theorem of Goddard and Thorn in string theory! A statement and proof of this theorem in an algebraic context can be found in E. Jurisich, Journal of Pure and Applied Algebra 126 (1998), 233-266).

Thus the graded components $\mathcal{M}_{(m, n)}$ of the Monster Lie algebra $\mathfrak{M}$ are as shown in Table 21.1. In this table $V_{n}$ is a vector space of dimension $c_{n}$ if $n \geq 1$ and $V_{-1}$ is a vector space of dimension 1.

Table 21.1 Graded components $\mathfrak{M}_{(m, n)}$ of the Monster Lie Algebra.

| $\cdots$ | $O$ | $O$ | $O$ | $O$ | $O$ | $V_{4}$ | $V_{8}$ | $V_{12}$ | $V_{16}$ | $\cdots$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
|  | $O$ | $O$ | $O$ | $O$ | $O$ | $V_{3}$ | $V_{6}$ | $V_{9}$ | $V_{12}$ |  |
|  | $O$ | $O$ | $O$ | $O$ | $O$ | $V_{2}$ | $V_{4}$ | $V_{6}$ | $V_{8}$ |  |
|  | $O$ | $O$ | $O$ | $V_{-1}$ | $O$ | $V_{1}$ | $V_{2}$ | $V_{3}$ | $V_{4}$ |  |
|  | $O$ | $O$ | $O$ | $O$ | $\mathbb{R}^{2}$ | $O$ | $O$ | $O$ | $O$ |  |
|  | $V_{4}$ | $V_{3}$ | $V_{2}$ | $V_{1}$ | $O$ | $V_{-1}$ | $O$ | $O$ | $O$ |  |
|  | $V_{8}$ | $V_{6}$ | $V_{4}$ | $V_{2}$ | $O$ | $O$ | $O$ | $O$ | $O$ |  |
|  | $V_{12}$ | $V_{9}$ | $V_{6}$ | $V_{3}$ | $O$ | $O$ | $O$ | $O$ | $O$ |  |
| $\cdots$ | $V_{16}$ | $V_{12}$ | $V_{8}$ | $V_{4}$ | $O$ | $O$ | $O$ | $O$ | $O$ | $\cdots$ |

We now consider the roots of the Monster Lie algebra $\mathfrak{M}$. Since $\mathfrak{M}=$ $L(\mathfrak{H}) / Z$ we recall from Section 21.2 that the root lattice of $\mathfrak{M}$ is defined to be the root lattice of $L(\mathfrak{H})$. The fundamental roots of $\mathfrak{M}$ are the $\alpha_{i}$ for $i \in I$. We have a homomorphism $Q \rightarrow H$ under which $\alpha_{i}$ maps to $h_{i}$. We pointed out in Section 21.2 that this homomorphism is not in general injective. In the Monster Lie algebra it is far from injective, as $\alpha_{i}, \alpha_{j}$ have the same image if and only if $i, j$ lie in the same block of $I$.

We have

$$
\begin{aligned}
& \left\langle\alpha_{0}, \alpha_{0}\right\rangle=2 \\
& \left\langle\alpha_{i}, \alpha_{i}\right\rangle=-2 m \quad \text { if } i \neq 0 \quad \text { and } i \in B_{m}
\end{aligned}
$$

Thus $\Pi_{\mathrm{re}}=\left\{\alpha_{0}\right\}$ and $\Pi_{\mathrm{im}}=\left\{\alpha_{i} ; i \neq 0\right\}$. Hence the Monster Lie algebra $\mathfrak{M}$ has just one real fundamental root and countably many imaginary fundamental roots.

The Weyl group $W$ of $\mathfrak{M}$ is generated by the fundamental reflections corresponding to the real fundamental roots. Thus $W=\left\langle s_{0}\right\rangle$, and so $W$ has order 2. Thus $\mathfrak{M}$ has an infinite number of fundamental roots while at the same time having a very small Weyl group isomorphic to the cyclic group of order 2.

Finally we shall consider Borcherds' denominator formula for the Monster Lie algebra $\mathfrak{M}$. This formula plays an important role in Borcherds’ proof of the Conway-Norton conjectures. We recall from Theorem 21.7 that Borcherds’ denominator formula is given by

$$
e(\rho) \prod_{\alpha \in \Phi^{+}}(1-e(-\alpha))^{m_{\alpha}}=\sum_{w \in W} \varepsilon(w) w\left(e(\rho) \sum_{\Psi}(-1)^{|\Psi|} e\left(-\sum \Psi\right)\right)
$$

where $\rho \in Q \otimes \mathbb{R}$ is any vector satisfying

$$
\left\langle\rho, \alpha_{i}\right\rangle=\frac{1}{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle \quad \text { for all } i \in I
$$

Now we have

$$
\bar{h}_{n}=\frac{(n+1) \bar{h}_{1}-(n-1) \bar{h}_{0}}{2}=b_{0}+n b_{1} .
$$

Thus we can identify a fundamental root $\alpha_{i}$ in the block $B_{n}$ with its image $b_{0}+n b_{1}$ in the Cartan subalgebra $\mathfrak{M}_{0}$ of $\mathfrak{M}$ provided we remember that there will be $c_{n}$ different such fundamental roots $\alpha_{i}$ with a given image $b_{0}+n b_{1}$.

We may take $\rho=\frac{\left(-\alpha_{0}-\alpha_{1}\right)}{2}$, since if $\alpha_{i} \in B_{n}$ we have

$$
\left\langle\rho, \alpha_{i}\right\rangle=\left\langle\frac{-\bar{h}_{0}-\bar{h}_{1}}{2}, b_{0}+n b_{1}\right\rangle=\left\langle-b_{0}, b_{0}+n b_{1}\right\rangle=n
$$

whereas

$$
\left\langle\alpha_{i}, \alpha_{i}\right\rangle=\left\langle b_{0}+n b_{1}, b_{0}+n b_{1}\right\rangle=-2 n
$$

Thus $\left\langle\rho, \alpha_{i}\right\rangle=\frac{1}{2}\left\langle\alpha_{i}, \alpha_{i}\right\rangle$.
Hence we shall use this vector $\rho$ in Borcherds' denominator formula. Using the natural homomorphisms

$$
\begin{aligned}
& Q \rightarrow H \rightarrow \mathfrak{M}_{0}=H / Z \\
& \alpha_{i} \rightarrow h_{i} \rightarrow \overline{h_{i}}
\end{aligned}
$$

we may interpret Borcherds' denominator formula in the integral group ring of $e\left(\mathfrak{M}_{0}\right)$ rather than the integral group ring of $e(Q)$. Bearing this in mind we define

$$
p=e^{b_{0}}, \quad q=e^{b_{1}}
$$

Then $e^{\rho}=p^{-1}$ and so the left-hand side of the denominator identity is

$$
p^{-1} \prod_{\substack{m>0 \\ n \in \mathbb{Z}}}\left(1-p^{m} q^{n}\right)^{c_{m n}}
$$

since the positive roots $\alpha \in \Phi^{+}$are the elements of $Q^{+}$which map to elements of $\mathfrak{M}_{0}$ of the form $m b_{0}+n b_{1}$ with $m>0$ and $n \in \mathbb{Z}$, and the number of $\alpha \in \Phi^{+}$ mapping to $m b_{0}+n b_{1}$ is $\operatorname{dim} \mathfrak{M}_{(m, n)}=c_{m n}$.

We now consider the right-hand side of the denominator identity. We recall that, for $\alpha \in Q, \varepsilon(\alpha)=(-1)^{k}$ where $\alpha$ is the sum of $k$ orthogonal imaginary fundamental roots and $\varepsilon(\alpha)=0$ otherwise. In the case of the Monster Lie algebra $\mathfrak{M}$ no two imaginary fundamental roots are orthogonal since

$$
\left\langle b_{0}+m b_{1}, b_{0}+n b_{1}\right\rangle=-(m+n)
$$

Thus the elements $\alpha \in Q$ contributing to $\sum \varepsilon(\alpha) e^{\alpha}$ are $\alpha=0$ with $\varepsilon(\alpha)=1$ and the imaginary simple roots in $Q$. These map to elements of form $b_{0}+n b_{1} \in$ $\mathfrak{M}_{0}$. There are $c_{n}$ such roots $\alpha \in Q$ mapping to $b_{0}+n b_{1}$ and they all give $\varepsilon(\alpha)=-1$. Thus

$$
\sum_{\alpha \in Q} \varepsilon(\alpha) e^{\alpha}=1-\sum_{n>0} c_{n} p q^{n}
$$

Now the Weyl group $W$ has order 2 and consists of the elements 1 and $s_{0}$. We have $s_{0}(p)=q$ and $s_{0}(q)=p$. Thus the right-hand side of the denominator identity is

$$
\begin{aligned}
& p^{-1}\left(1-\sum_{n>0} c_{n} p q^{n}\right)-q^{-1}\left(1-\sum_{n>0} c_{n} q p^{n}\right) \\
& \quad=\left(p^{-1}+\sum_{n>0} c_{n} p^{n}\right)-\left(q^{-1}+\sum_{n>0} c_{n} q^{n}\right) \\
& \quad=j(p)-j(q)
\end{aligned}
$$

Thus we have obtained the following result.

Theorem 21.10 Borcherds' denominator identity for the Monster Lie algebra $\mathfrak{M}$ asserts that:

$$
p^{-1} \prod_{\substack{m>0 \\ n \in \mathbb{Z}}}\left(1-p^{m} q^{n}\right)^{c_{m n}}=j(p)-j(q)
$$

where $c_{n}$ is the coefficient of $q^{n}$ in the modular function $j$.
In fact this identity was proved by Borcherds from first principles and used subsequently to prove that the fundamental roots of $\mathfrak{M}$ map to the elements

$$
b_{0}-b_{1}, b_{0}+b_{1}, b_{0}+2 b_{1}, b_{0}+3 b_{1}, \ldots
$$

of $\mathfrak{M}_{0}$.
Further information about results stated without proof in this chapter can be found in the papers of R. Borcherds 'Generalised Kac-Moody algebras', Journal of Algebra 115 (1988), 501-512, and 'Monstrous moonshine and monstrous Lie superalgebras’, Inventiones Mathematiae 109 (1992), 405-444.

## Appendix

## Summary pages - explanation

There follow a number of summary pages, one for each Lie algebra of finite or affine type, giving basic properties of the Lie algebra in question. The information given differs to some extent between the Lie algebras of finite type and those of affine type.

In the case of the algebras of finite type we give the name of the algebra, the Dynkin diagram with the labelling we have chosen for its vertices, the Cartan matrix, the dimension of the Lie algebra, its Coxeter number, the order of its Weyl group $W$ and the degrees of the basic polynomial invariants of $W$. We also give information about its root system. The roots are most conveniently described in terms of a basis $\beta_{1}, \ldots, \beta_{m}$ of mutually orthogonal basis vectors all of the same length. In several cases it is convenient to choose $m$ greater than the rank $l$ of the Lie algebra, so that the root system lies in a proper subspace of the vector space spanned by $\beta_{1}, \ldots, \beta_{m}$. In the cases when there are roots of two different lengths the long roots and short roots are both described. The extended Dynkin diagram is given and the root lattice described in terms of the above orthogonal basis. The fundamental weights are given, as is the index of the root lattice in the weight lattice. Finally the standard invariant forms on $H_{\mathbb{R}}$ and $H_{\mathbb{R}}^{*}$ are described, and the constant is given which converts the standard invariant form on $H_{\mathbb{R}}$ into the Killing form.

The labelling given here for the vertices of the Dynkin diagrams of types $E_{6}$ and $E_{7}$ differs from that used in Chapter 8, where it was convenient to describe the root systems of type $E_{6}, E_{7}$ or $E_{8}$ together in Section 8.7.

In the case of the Lie algebras of affine type we have given two names for each algebra which we have called the Dynkin name and the Kac name. The Dynkin name describes the Dynkin diagram of the algebra whereas the Kac name, introduced at the end of Chapter 18, indicates whether the Lie algebra is
of untwisted or twisted type, and in the case of those of twisted type indicates the type of the untwisted affine algebra from which it is obtained, together with the order of the automorphism of which it is the fixed point subalgebra. This Kac notation is entirely consistent with the notation normally used to describe the twisted Chevalley groups.

The Dynkin diagram with chosen labelling is given, together with the generalised Cartan matrix and the integers $a_{0}, a_{1}, \ldots, a_{l}$ and $c_{0}, c_{1}, \ldots, c_{l}$. The central element $c$, the basic imaginary root $\delta$, and the elements $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$ and $h_{\theta} \in H_{\mathbb{R}}^{0}$ which play an important role in the theory of affine algebras are written down explicitly. The Coxeter number and dual Coxeter number are given.

The type of the finite dimensional Lie algebra $L^{0}$ obtained by removing vertex 0 from the Dynkin diagram is given. The root system $\Phi$ is described in terms of the root system $\Phi^{0}$ of $L^{0}$. The real and imaginary roots are given separately and the multiplicities of the imaginary roots are given. (The real roots all have multiplicity 1.) In order to clarify the action of the affine Weyl group we describe the lattices $M \subset H_{\mathbb{R}}^{0}$ and $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$ which give rise to the translations in the affine Weyl group. We also describe the fundamental alcoves $A \subset H_{\mathbb{R}}^{0}$ and $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$ whose closures give fundamental regions for the action of the affine Weyl group. Then we describe the fundamental weights in terms of the fundamental weights of $L^{0}$, and the standard invariant forms on $H$ and on $H^{*}$.

For the affine algebras of types $\tilde{C}_{l}^{\prime}, l \geq 2$, and $\tilde{A}_{1}^{\prime}$ we have given two different descriptions, corresponding to two choices of the vertex of the Dynkin diagram labelled by 0 . (The algebra $\tilde{A}_{1}^{\prime}$ behaves just like $\tilde{C}_{l}^{\prime}$ when $l=1$.) Both descriptions are useful, as is shown in Section 18.4. The first description is the conventional description in which the associated finite dimensional algebra has type $C_{l}$, and which is discussed in Chapter 17. The second description is the one used to obtain the realisation of $\tilde{C}_{l}^{\prime}$ as ${ }^{2} \tilde{A}_{2 l}$ in Section 18.4. Here the associated finite dimensional algebra has type $B_{l}$. A word of caution is necessary in deriving the results appearing in the second description. In these cases we have $c_{0}=2$. Thus we cannot apply results from Chapter 17 uncritically to these cases, since $c_{0}=1$ is assumed in Chapter 17. Instead we have the following situation.

$$
\theta=\delta-a_{0} \alpha_{0}=\sum_{i=1}^{l} a_{i} \alpha_{i}
$$

satisfies $\langle\theta, \theta\rangle=2 a_{0} c_{0}$. We also have

$$
h_{\theta}=\frac{1}{a_{0} c_{0}}\left(c-c_{0} h_{0}\right)=\frac{1}{a_{0} c_{0}} \sum_{i=1}^{l} c_{i} h_{i} .
$$

Under the natural bijection $H \leftrightarrow H^{*}$ we have $h_{i} \leftrightarrow a_{i} c_{i}^{-1} \alpha_{i}, h_{\theta} \leftrightarrow a_{0}^{-1} c_{0}^{-1} \theta$, $d \leftrightarrow a_{0} c_{0}^{-1} \gamma$. In addition we have

$$
\left\langle h_{0}, d\right\rangle=a_{0} c_{0}^{-1},\left\langle\alpha_{0}, \gamma\right\rangle=a_{0}^{-1} c_{0}
$$

The lattices $M, M^{*}$ are given as follows in these cases. $M$ is the lattice generated by $w\left(h_{\theta}\right)$ for all $w \in W^{0}$, and $M^{*}$ is the lattice generated by $w\left(a_{0}^{-1} c_{0}^{-1} \theta\right)$ for all $w \in W^{0}$. The alcove $A$ is bounded by the affine hyperplane $\theta(h)=1$ and the alcove $A^{*}$ is bounded by the affine hyperplane $\lambda\left(h_{\theta}\right)=\frac{1}{a_{0} c_{0}}$.

In fact in this second description $\theta$ turns out to be $2 \theta_{\mathrm{s}}$ where $\theta_{\mathrm{s}}$ is the highest short root, and $h_{\theta}$ is $\frac{1}{2} h_{\theta_{\mathrm{s}}}$.

NAME $A_{l}$
Dynkin diagram with labelling.


Cartan matrix


Dimension. $\quad \operatorname{dim} L=l(l+2)$.
Coxeter number. $\quad h=l+1$.
Order of the Weyl group. $\quad|W|=(l+1)$ !
Degrees of the basic polynomial invariants of $W$.

$$
\left\{d_{1}, d_{2}, \ldots, d_{l}\right\}=\{2,3, \ldots, l+1\}
$$

Number of roots. $\quad|\Phi|=l(l+1)$.
The fundamental roots in terms of an orthogonal basis.

$$
\alpha_{1}=\beta_{1}-\beta_{2}, \quad \alpha_{2}=\beta_{2}-\beta_{3}, \quad \ldots, \quad \alpha_{l}=\beta_{l}-\beta_{l+1} .
$$

The root system.

$$
\Phi=\left\{\beta_{i}-\beta_{j} ; \quad i, j=1, \ldots, l+1, \quad i \neq j\right\}
$$

The highest root. $\quad \theta=\beta_{1}-\beta_{l+1}$.

The extended Dynkin diagram, for $l \geq 2$.


The root lattice $\quad Q=\sum \mathbb{Z} \alpha_{i}$.

$$
Q=\left\{\sum_{i=1}^{l+1} \xi_{i} \beta_{i} ; \quad \xi_{i} \in \mathbb{Z}, \sum \xi_{i}=0\right\}
$$

The fundamental weights.

$$
\omega_{i}=\frac{1}{l+1}\left((l+1-i)\left(\beta_{1}+\cdots+\beta_{i}\right)-i\left(\beta_{i+1}+\cdots+\beta_{l+1}\right)\right) \quad i=1, \ldots, l
$$

The index of the root lattice in the weight lattice. $\quad|X: Q|=l+1 . \quad X / Q$ is cyclic.

The standard invariant form on $H_{\mathbb{R}}$.

$$
\left\langle h_{i}, h_{j}\right\rangle=A_{i j}
$$

The standard invariant form on $H_{\mathbb{R}}^{*}$.

$$
\left\langle\alpha_{i}, \alpha_{j}\right\rangle=A_{i j}
$$

The Killing form on $H_{\mathbb{R}}$.

$$
\langle x, y\rangle_{K}=\frac{1}{b}\langle x, y\rangle
$$

where $b=2(l+1)$.

NAME $B_{l}$
Dynkin diagram with labelling.


Cartan matrix


Dimension. $\quad \operatorname{dim} L=l(2 l+1)$.
Coxeter number. $\quad h=2 l$.
Order of the Weyl group. $\quad|W|=2^{l} \cdot l$ !
Degrees of the basic polynomial invariants of $W$.

$$
\left\{d_{1}, d_{2}, \ldots, d_{l}\right\}=\{2,4, \ldots, 2 l\}
$$

Number of roots. $\quad|\Phi|=2 l^{2}$.
The fundamental roots in terms of an orthogonal basis.

$$
\alpha_{1}=\beta_{1}-\beta_{2}, \quad \alpha_{2}=\beta_{2}-\beta_{3}, \quad \ldots, \quad \alpha_{l-1}=\beta_{l-1}-\beta_{l}, \quad \alpha_{l}=\beta_{l}
$$

The root system. $\quad \Phi=\Phi_{1} \cup \Phi_{\text {s }}$ where

$$
\begin{aligned}
& \Phi_{1}=\left\{ \pm \beta_{i} \pm \beta_{j} ; \quad i, \quad j=1, \ldots, l, \quad i \neq j\right\} \\
& \Phi_{\mathrm{s}}=\left\{ \pm \beta_{i} ; i=1, \ldots, l\right\}
\end{aligned}
$$

The highest root. $\quad \theta_{1}=\beta_{1}+\beta_{2}$.
The highest short root. $\quad \theta_{\mathrm{s}}=\beta_{1}$.

The extended Dynkin diagram.


The root lattice $\quad Q=\sum \mathbb{Z} \alpha_{i}$.

$$
Q=\left\{\sum_{i=1}^{l} \xi_{i} \beta_{i} ; \quad \xi_{i} \in \mathbb{Z}\right\}
$$

The fundamental weights.

$$
\begin{aligned}
& \omega_{i}=\beta_{1}+\cdots+\beta_{i} \quad i=1, \ldots, l-1 \\
& \omega_{l}=\frac{1}{2}\left(\beta_{1}+\cdots+\beta_{l}\right)
\end{aligned}
$$

The index of the root lattice in the weight lattice. $\quad|X: Q|=2$.
The symmetrising matrix $D=\operatorname{diag}\left(d_{i}\right)$.

$$
d_{i}=1 \quad i=1, \ldots, l-1 \quad d_{l}=2 .
$$

The standard invariant form on $H_{\mathbb{R}}$.

$$
\left\langle h_{i}, h_{j}\right\rangle=A_{i j} d_{j} .
$$

The standard invariant form on $H_{\mathbb{R}}^{*}$.

$$
\left\langle\alpha_{i}, \alpha_{j}\right\rangle=d_{i}^{-1} A_{i j}
$$

The Killing form on $H_{\mathbb{R}}$.

$$
\langle x, y\rangle_{K}=\frac{1}{b}\langle x, y\rangle
$$

where $b=4 l-2$.

NAME $C_{l}$
Dynkin diagram with labelling.


Cartan matrix


Dimension. $\quad \operatorname{dim} L=l(2 l+1)$.
Coxeter number. $\quad h=2 l$.
Order of the Weyl group. $\quad|W|=2^{l} \cdot l$ !
Degrees of the basic polynomial invariants of $W$.

$$
\left\{d_{1}, d_{2}, \ldots, d_{l}\right\}=\{2,4, \ldots, 2 l\}
$$

Number of roots. $\quad|\Phi|=2 l^{2}$.
The fundamental roots in terms of an orthogonal basis.

$$
\alpha_{1}=\beta_{1}-\beta_{2}, \quad \alpha_{2}=\beta_{2}-\beta_{3}, \quad \ldots, \quad \alpha_{l-1}=\beta_{l-1}-\beta_{l}, \quad \alpha_{l}=2 \beta_{l} .
$$

The root system. $\quad \Phi=\Phi_{1} \cup \Phi_{\text {s }}$ where

$$
\begin{aligned}
& \Phi_{1}=\left\{ \pm 2 \beta_{i} ; i=1, \ldots, l\right\} \\
& \Phi_{\mathrm{s}}=\left\{ \pm \beta_{i} \pm \beta_{j} ; \quad i, j=1, \ldots, l \quad i \neq j\right\}
\end{aligned}
$$

The highest root. $\quad \theta_{1}=2 \beta_{1}$.
The highest short root. $\quad \theta_{\mathrm{s}}=\beta_{1}+\beta_{2}$.

The extended Dynkin diagram.


The root lattice $\quad Q=\sum \mathbb{Z} \alpha_{i}$.

$$
Q=\left\{\sum_{i=1}^{l} \xi_{i} \beta_{i} ; \xi_{i} \in \mathbb{Z}, \sum \xi_{i} \text { even }\right\}
$$

The fundamental weights

$$
\omega_{i}=\beta_{1}+\cdots+\beta_{i} \quad i=1, \ldots, l .
$$

The index of the root lattice in the weight lattice. $\quad|X: Q|=2$.
The symmetrising matrix $D=\operatorname{diag}\left(d_{i}\right)$.

$$
d_{i}=2 \quad i=1, \ldots, l-1 \quad d_{l}=1
$$

The standard invariant form on $H_{\mathbb{R}}$.

$$
\left\langle h_{i}, h_{j}\right\rangle=A_{i j} d_{j}
$$

The standard invariant form on $H_{\mathbb{R}}^{*}$.

$$
\left\langle\alpha_{i}, \alpha_{j}\right\rangle=d_{i}^{-1} A_{i j}
$$

The Killing form on $H_{\mathbb{R}}$.

$$
\langle x, y\rangle_{K}=\frac{1}{b}\langle x, y\rangle
$$

where $b=2(l+1)$.

NAME $D_{l}$
Dynkin diagram with labelling.


Cartan matrix


Dimension. $\quad \operatorname{dim} L=l(2 l-1)$.
Coxeter number. $\quad h=2 l-2$.
Order of the Weyl group. $\quad|W|=2^{l-1} l$ !
Degrees of the basic polynomial invariants of $W$.

$$
\left\{d_{1}, d_{2}, \ldots, d_{l}\right\}=\{2,4, \ldots, 2 l-2, l\}
$$

Number of roots. $\quad|\Phi|=2 l(l-1)$.
The fundamental roots in terms of an orthogonal basis.

$$
\begin{aligned}
\alpha_{1} & =\beta_{1}-\beta_{2}, \quad \alpha_{2}=\beta_{2}-\beta_{3}, \quad \ldots, \quad \alpha_{l-2}=\beta_{l-2}-\beta_{l-1}, \\
\alpha_{l-1} & =\beta_{l-1}-\beta_{l}, \quad \alpha_{l}=\beta_{l-1}+\beta_{l} .
\end{aligned}
$$

The root system.

$$
\Phi=\left\{ \pm \beta_{i} \pm \beta_{j} ; \quad i, j=1, \ldots, l \quad i \neq j\right\}
$$

The highest root. $\quad \theta=\beta_{1}+\beta_{2}$.
The extended Dynkin diagram.


$l=4$

The root lattice $Q=\sum \mathbb{Z} \alpha_{i}$.

$$
Q=\left\{\sum_{i=1}^{l} \xi_{i} \beta_{i} ; \quad \xi_{i} \in \mathbb{Z}, \quad \sum \xi_{i} \text { even }\right\}
$$

The fundamental weights.

$$
\begin{aligned}
\omega_{i} & =\beta_{1}+\cdots+\beta_{i} \quad i=1, \ldots, l-2 \\
\omega_{l-1} & =\frac{1}{2}\left(\beta_{1}+\cdots+\beta_{l-2}+\beta_{l-1}-\beta_{l}\right) \\
\omega_{l} & =\frac{1}{2}\left(\beta_{1}+\cdots+\beta_{l-2}+\beta_{l-1}+\beta_{l}\right)
\end{aligned}
$$

The index of the root lattice in the weight lattice. $\quad|X: Q|=4$.
$X / Q$ is cyclic if $l$ is odd and non-cyclic if $l$ is even.
The standard invariant form on $H_{\mathbb{R}}$.

$$
\left\langle h_{i}, h_{j}\right\rangle=A_{i j}
$$

The standard invariant form on $H_{\mathbb{R}}^{*}$.

$$
\left\langle\alpha_{i}, \alpha_{j}\right\rangle=A_{i j}
$$

The Killing form on $H_{\mathbb{R}}$.

$$
\langle x, y\rangle_{K}=\frac{1}{b}\langle x, y\rangle
$$

where $b=4(l-1)$.

NAME $E_{6}$
Dynkin diagram with labelling.


Cartan matrix
$\left.\begin{array}{l}1 \\ 1 \\ 2 \\ 3 \\ 4 \\ 4 \\ 5 \\ 6 \\ 6 \\ 6\end{array} \begin{array}{rrrrrr}1 & 2 & 3 & 4 & 5 & 6 \\ 0 & -1 & 0 & 0 & 0 & 0 \\ 0 & 0 & -1 & 0 & 0 & 0 \\ 0 & 0 & -1 & 2 & 0 & 0 \\ 0 & 0 & 0 & 0 & -1 & 0 \\ 0\end{array}\right)$

Dimension. $\quad \operatorname{dim} L=78$.
Coxeter number. $\quad h=12$.
Order of the Weyl group. $|W|=2^{7} \cdot 3^{4} \cdot 5$
Degrees of the basic polynomial invariants of $W$.

$$
\left\{d_{1}, d_{2}, \ldots, d_{6}\right\}=\{2,5,6,8,9,12\} .
$$

Number of roots. $\quad|\Phi|=72$.
The fundamental roots in terms of an orthogonal basis.
$\beta_{1}, \beta_{2}, \beta_{3}, \beta_{4}, \beta_{5}, \beta_{6}, \beta_{7}, \beta_{8} \quad$ orthogonal basis.

$$
\begin{array}{lll}
\alpha_{1}=\beta_{1}-\beta_{2}, & \alpha_{2}=\beta_{2}-\beta_{3}, \quad \alpha_{3}=\beta_{3}-\beta_{4}, \quad \alpha_{4}=\beta_{4}-\beta_{5}, \\
\alpha_{5}=\beta_{4}+\beta_{5}, & \alpha_{6}=-\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}+\beta_{7}+\beta_{8}\right) .
\end{array}
$$

The root system.

$$
\begin{aligned}
\Phi= & \left\{ \pm \beta_{i} \pm \beta_{j} ; \quad i, j=1,2,3,4,5 \quad i \neq j\right\} \\
& \cup\left\{\frac{1}{2} \sum_{i=1}^{8} \varepsilon_{i} \beta_{i} ; \quad \varepsilon_{i}= \pm 1, \prod_{i=1}^{8} \varepsilon_{i}=1, \quad \varepsilon_{6}=\varepsilon_{7}=\varepsilon_{8}\right\}
\end{aligned}
$$

The highest root.

$$
\theta=\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}-\beta_{5}-\beta_{6}-\beta_{7}-\beta_{8}\right)
$$

The extended Dynkin diagram.


The root lattice $Q=\sum \mathbb{Z} \alpha_{i}$.
$Q=\left\{\sum_{i=1}^{8} \xi_{i} \beta_{i} ; 2 \xi_{i} \in \mathbb{Z}, \xi_{i}-\xi_{j} \in \mathbb{Z}, \sum_{i=1}^{8} \xi_{i} \in 2 \mathbb{Z}, i, j=1, \ldots, 8 \xi_{6}=\xi_{7}=\xi_{8}\right\}$.
The fundamental weights.

$$
\begin{aligned}
& \omega_{1}=\beta_{1}-\frac{1}{3}\left(\beta_{6}+\beta_{7}+\beta_{8}\right) \\
& \omega_{2}=\beta_{1}+\beta_{2}-\frac{2}{3}\left(\beta_{6}+\beta_{7}+\beta_{8}\right) \\
& \omega_{3}=\beta_{1}+\beta_{2}+\beta_{3}-\left(\beta_{6}+\beta_{7}+\beta_{8}\right) \\
& \omega_{4}=\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}-\beta_{5}-\beta_{6}-\beta_{7}-\beta_{8}\right) \\
& \omega_{5}=\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}\right)-\frac{5}{6}\left(\beta_{6}+\beta_{7}+\beta_{8}\right) \\
& \omega_{6}=-\frac{2}{3}\left(\beta_{6}+\beta_{7}+\beta_{8}\right)
\end{aligned}
$$

The index of the root lattice in the weight lattice. $\quad|X: Q|=3$.
The standard invariant form on $H_{\mathbb{R}}$.

$$
\left\langle h_{i}, h_{j}\right\rangle=A_{i j}
$$

The standard invariant form on $H_{\mathbb{R}}^{*}$.

$$
\left\langle\alpha_{i}, \alpha_{j}\right\rangle=A_{i j}
$$

The Killing form on $H_{\mathbb{R}}$.

$$
\langle x, y\rangle_{K}=\frac{1}{b}\langle x, y\rangle
$$

where $b=24$.

NAME $E_{7}$
Dynkin diagram with labelling.


Cartan matrix
1
2
2
4
4
5
6
7
7 $\left(\begin{array}{rrrrrrr}1 & 2 & 3 & 4 & 5 & 6 & 7 \\ 2 & -1 & 0 & 0 & 0 & 0 & 0 \\ -1 & 2 & -1 & 0 & 0 & 0 & 0 \\ 0 & -1 & 2 & -1 & 0 & 0 & 0 \\ 0 & 0 & -1 & 2 & -1 & -1 & 0 \\ 0 & 0 & 0 & -1 & 2 & 0 & 0 \\ 0 & 0 & 0 & -1 & 0 & 2 & -1 \\ 0 & 0 & 0 & 0 & -1 & 2\end{array}\right)$

Dimension. $\operatorname{dim} L=133$.
Coxeter number. $\quad h=18$.
Order of the Weyl group. $\quad|W|=2^{10} \cdot 3^{4} \cdot 5 \cdot 7$.
Degrees of the basic polynomial invariants of $W$.

$$
\left\{d_{1}, d_{2}, \ldots, d_{7}\right\}=\{2,6,8,10,12,14,18\} .
$$

Number of roots. $\quad|\Phi|=126$.
The fundamental roots in terms of an orthogonal basis.
$\beta_{1}, \beta_{2}, \beta_{3}, \beta_{4}, \beta_{5}, \beta_{6}, \beta_{7}, \beta_{8} \quad$ orthogonal basis.

$$
\begin{array}{ll}
\alpha_{1}=\beta_{1}-\beta_{2}, & \alpha_{2}=\beta_{2}-\beta_{3}, \quad \alpha_{3}=\beta_{3}-\beta_{4}, \quad \alpha_{4}=\beta_{4}-\beta_{5}, \quad \alpha_{5}=\beta_{5}-\beta_{6}, \\
\alpha_{6}=\beta_{5}+\beta_{6}, & \alpha_{7}=-\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}+\beta_{7}+\beta_{8}\right) .
\end{array}
$$

The root system.

$$
\begin{aligned}
\Phi= & \left\{ \pm \beta_{i} \pm \beta_{j} ; \quad i, j=1,2,3,4,5,6 \quad i \neq j\right\} \\
& \cup\left\{ \pm\left(\beta_{7}+\beta_{8}\right)\right\} \\
& \cup\left\{\frac{1}{2} \sum_{i=1}^{8} \varepsilon_{i} \beta_{i} ; \quad \varepsilon_{i}= \pm 1, \prod_{i=1}^{8} \varepsilon_{i}=1, \varepsilon_{7}=\varepsilon_{8}\right\}
\end{aligned}
$$

The highest root. $\quad \theta=-\beta_{7}-\beta_{8}$.
The extended Dynkin diagram


The root lattice $Q=\sum \mathbb{Z} \alpha_{i}$.

$$
Q=\left\{\sum_{i=1}^{8} \xi_{i} \beta_{i} ; 2 \xi_{i} \in \mathbb{Z}, \quad \xi_{i}-\xi_{j} \in \mathbb{Z}, \sum_{i=1}^{8} \xi_{i} \in 2 \mathbb{Z}, i, j=1, \ldots, 8 \quad \xi_{7}=\xi_{8}\right\}
$$

The fundamental weights.

$$
\begin{aligned}
& \omega_{1}=\beta_{1}-\frac{1}{2}\left(\beta_{7}+\beta_{8}\right) \\
& \omega_{2}=\beta_{1}+\beta_{2}-\left(\beta_{7}+\beta_{8}\right) \\
& \omega_{3}=\beta_{1}+\beta_{2}+\beta_{3}-\frac{3}{2}\left(\beta_{7}+\beta_{8}\right) \\
& \omega_{4}=\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}-2\left(\beta_{7}+\beta_{8}\right) \\
& \omega_{5}=\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}-\beta_{6}\right)-\left(\beta_{7}+\beta_{8}\right) \\
& \omega_{6}=\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}\right)-\frac{3}{2}\left(\beta_{7}+\beta_{8}\right) \\
& \omega_{7}=-\left(\beta_{7}+\beta_{8}\right)
\end{aligned}
$$

The index of the root lattice in the weight lattice. $\quad|X: Q|=2$.
The standard invariant form on $H_{\mathbb{R}}$.

$$
\left\langle h_{i}, h_{j}\right\rangle=A_{i j}
$$

The standard invariant form on $H_{\mathbb{R}}^{*}$.

$$
\left\langle\alpha_{i}, \alpha_{j}\right\rangle=A_{i j}
$$

The Killing form on $H_{\mathbb{R}}$.

$$
\langle x, y\rangle_{K}=\frac{1}{b}\langle x, y\rangle
$$

where $b=36$.

NAME $E_{8}$
Dynkin diagram with labelling.


Cartan matrix
$\left.\begin{array}{l}1 \\ 2 \\ 2 \\ 2 \\ 4 \\ 5 \\ 6 \\ 6 \\ 7 \\ 8 \\ 8 \\ -1\end{array} \begin{array}{rrrrrrrr}1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 \\ 0 & -1 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & -1 & 0 & 0 & 0 & 0 & 0 \\ 0 & 2 & -1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & -1 & 2 & -1 & 0 & 0 \\ 0 & 0 & -1 & -1 & 0 \\ 0 & 0 & 0 & -1 & 2 & 0 & 0 \\ 0 & 0 & 0 & -1 & 0 & 2 & -1 \\ 0 & 0 & 0 & 0 & -1 & 2\end{array}\right)$

Dimension. $\operatorname{dim} L=248$.
Coxeter number. $\quad h=30$.
Order of the Weyl group. $\quad|W|=2^{14} \cdot 3^{5} \cdot 5^{2} \cdot 7$.
Degrees of the basic polynomial invariants of $W$.

$$
\left\{d_{1}, d_{2}, \ldots, d_{8}\right\}=\{2,8,12,14,18,20,24,30\} .
$$

Number of roots. $\quad|\Phi|=240$.
The fundamental roots in terms of an orthogonal basis.

$$
\begin{array}{lll}
\alpha_{1}=\beta_{1}-\beta_{2}, & \alpha_{2}=\beta_{2}-\beta_{3}, \quad \alpha_{3}=\beta_{3}-\beta_{4}, \quad \alpha_{4}=\beta_{4}-\beta_{5}, \\
\alpha_{5}=\beta_{5}-\beta_{6}, & \alpha_{6}=\beta_{6}-\beta_{7}, \quad \alpha_{7}=\beta_{6}+\beta_{7}, \\
\alpha_{8}=-\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}+\beta_{7}+\beta_{8}\right) .
\end{array}
$$

The root system.

$$
\begin{aligned}
\Phi= & \left\{ \pm \beta_{i} \pm \beta_{j} ; \quad i, j=1,2,3,4,5,6,7,8 \quad i \neq j\right\} \\
& \cup\left\{\frac{1}{2} \sum_{i=1}^{8} \varepsilon_{i} \beta_{i} ; \quad \varepsilon_{i}= \pm 1, \prod_{i=1}^{8} \varepsilon_{i}=1\right\}
\end{aligned}
$$

The highest root. $\quad \theta=\beta_{1}-\beta_{8}$.
The extended Dynkin diagram


The root lattice $Q=\sum \mathbb{Z} \alpha_{i}$.

$$
Q=\left\{\sum_{i=1}^{8} \xi_{i} \beta_{i} ; 2 \xi_{i} \in \mathbb{Z}, \quad \xi_{i}-\xi_{j} \in \mathbb{Z}, \quad \sum_{i=1}^{8} \xi_{i} \in 2 \mathbb{Z} \quad i, j=1, \ldots, 8\right\}
$$

The fundamental weights.

$$
\begin{aligned}
& \omega_{1}=\beta_{1}-\beta_{8} \\
& \omega_{2}=\beta_{1}+\beta_{2}-2 \beta_{8} \\
& \omega_{3}=\beta_{1}+\beta_{2}+\beta_{3}-3 \beta_{8} \\
& \omega_{4}=\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}-4 \beta_{8} \\
& \omega_{5}=\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}-5 \beta_{8} \\
& \omega_{6}=\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}-\beta_{7}\right)-\frac{5}{2} \beta_{8} \\
& \omega_{7}=\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+\beta_{4}+\beta_{5}+\beta_{6}+\beta_{7}\right)-\frac{7}{2} \beta_{8} \\
& \omega_{8}=-2 \beta_{8}
\end{aligned}
$$

The index of the root lattice in the weight lattice. $\quad|X: Q|=1$.
The standard invariant form on $H_{\mathbb{R}}$.

$$
\left\langle h_{i}, h_{j}\right\rangle=A_{i j}
$$

The standard invariant form on $H_{\mathbb{R}}^{*}$.

$$
\left\langle\alpha_{i}, \alpha_{j}\right\rangle=A_{i j}
$$

The Killing form on $H_{\mathbb{R}}$.

$$
\langle x, y\rangle_{K}=\frac{1}{b}\langle x, y\rangle
$$

where $b=60$.

NAME $F_{4}$
Dynkin diagram with labelling.


Cartan matrix
$\left.\begin{array}{c}1 \\ 2 \\ 3 \\ 4 \\ 0 \\ -1\end{array} \begin{array}{rrrr}1 & 2 & 3 & 4 \\ 0 & -1 & 0 & 0 \\ 0 & -1 & 0 \\ 0 & -1 & -1 \\ 2\end{array}\right)$

Dimension. $\operatorname{dim} L=52$.
Coxeter number. $\quad h=12$.
Order of the Weyl group. $\quad|W|=2^{7} \cdot 3^{2}$.
Degrees of the basic polynomial invariants of $W$.

$$
\left\{d_{1}, d_{2}, d_{3}, d_{4}\right\}=\{2,6,8,12\}
$$

Number of roots. $\quad|\Phi|=48$.
The fundamental roots in terms of an orthogonal basis.

$$
\alpha_{1}=\beta_{1}-\beta_{2}, \quad \alpha_{2}=\beta_{2}-\beta_{3}, \quad \alpha_{3}=\beta_{3}, \quad \alpha_{4}=\frac{1}{2}\left(-\beta_{1}-\beta_{2}-\beta_{3}+\beta_{4}\right) .
$$

The root system. $\quad \Phi=\Phi_{1} \cup \Phi_{\mathrm{s}}$ where

$$
\begin{aligned}
& \Phi_{1}=\left\{ \pm \beta_{i} \pm \beta_{j} ; \quad i, j=1,2,3,4 \quad i \neq j\right\} \\
& \Phi_{\mathrm{s}}=\left\{ \pm \beta_{i} \quad i=1,2,3,4\right\} \cup\left\{\frac{1}{2} \sum_{i=1}^{4} \varepsilon_{i} \beta_{i} ; \quad \varepsilon_{i}= \pm 1\right\}
\end{aligned}
$$

The highest root. $\quad \theta_{1}=\beta_{1}+\beta_{4}$.
The highest short root. $\quad \theta_{\mathrm{s}}=\beta_{4}$.
The extended Dynkin diagram.


The root lattice $Q=\sum \mathbb{Z} \alpha_{i}$.

$$
Q=\left\{\sum_{i=1}^{4} \xi_{i} \beta_{i} ; 2 \xi_{i} \in \mathbb{Z}, \quad \xi_{i}-\xi_{j} \in \mathbb{Z} i, j=1,2,3,4\right\}
$$

The fundamental weights.

$$
\begin{aligned}
& \omega_{1}=\beta_{1}+\beta_{4} \\
& \omega_{2}=\beta_{1}+\beta_{2}+2 \beta_{4} \\
& \omega_{3}=\frac{1}{2}\left(\beta_{1}+\beta_{2}+\beta_{3}+3 \beta_{4}\right) \\
& \omega_{4}=\beta_{4} .
\end{aligned}
$$

The index of the root lattice in the weight lattice. $\quad|X: Q|=1$.
The symmetrising matrix $D=\operatorname{diag}\left(d_{i}\right)$.

$$
d_{1}=1, \quad d_{2}=1, \quad d_{3}=2, \quad d_{4}=2
$$

The standard invariant form on $H_{\mathbb{R}}$.

$$
\left\langle h_{i}, h_{j}\right\rangle=A_{i j} d_{j}
$$

The standard invariant form on $H_{\mathbb{R}}^{*}$.

$$
\left\langle\alpha_{i}, \alpha_{j}\right\rangle=d_{i}^{-1} A_{i j}
$$

The Killing form on $H_{\mathbb{R}}$.

$$
\langle x, y\rangle_{K}=\frac{1}{b}\langle x, y\rangle
$$

where $b=18$.

## NAME $G_{2}$

Dynkin diagram with labelling.


Cartan matrix

$$
\begin{gathered}
1 \\
1\left(\begin{array}{cc}
2 \\
2 & -1 \\
-3
\end{array}\right)
\end{gathered}
$$

Dimension. $\operatorname{dim} L=14$.
Coxeter number. $\quad h=6$.
Order of the Weyl group. $\quad|W|=12$.
Degrees of the basic polynomial invariants of $W$.

$$
\left\{d_{1}, d_{2}\right\}=\{2,6\}
$$

Number of roots. $\quad|\Phi|=12$.
The fundamental roots in terms of an orthogonal basis.

$$
\begin{aligned}
& \beta_{1}, \beta_{2}, \beta_{3} \text { orthogonal basis. } \\
& \alpha_{1}=-2 \beta_{1}+\beta_{2}+\beta_{3}, \quad \alpha_{2}=\beta_{1}-\beta_{2}
\end{aligned}
$$

The root system. $\quad \Phi=\Phi_{1} \cup \Phi_{\text {s }}$ where

$$
\begin{aligned}
& \Phi_{1}=\left\{ \pm\left(-2 \beta_{1}+\beta_{2}+\beta_{3}\right), \quad \pm\left(\beta_{1}-2 \beta_{2}+\beta_{3}\right), \quad \pm\left(\beta_{1}+\beta_{2}-2 \beta_{3}\right)\right\} \\
& \Phi_{\mathrm{s}}=\left\{ \pm\left(\beta_{1}-\beta_{2}\right), \quad \pm\left(\beta_{2}-\beta_{3}\right), \quad \pm\left(\beta_{1}-\beta_{3}\right)\right\}
\end{aligned}
$$

The highest root. $\quad \theta_{1}=-\beta_{1}-\beta_{2}+2 \beta_{3}$.
The highest short root. $\quad \theta_{\mathrm{s}}=-\beta_{2}+\beta_{3}$.
The extended Dynkin diagram.


The root lattice $Q=\sum \mathbb{Z} \alpha_{i}$.

$$
Q=\left\{\sum_{i=1}^{3} \xi_{i} \beta_{i} ; \quad \xi_{i} \in \mathbb{Z}, \quad \xi_{1}+\xi_{2}+\xi_{3}=0\right\}
$$

The fundamental weights.

$$
\begin{aligned}
& \omega_{1}=-\beta_{1}-\beta_{2}+2 \beta_{3} \\
& \omega_{2}=-\beta_{2}+\beta_{3} .
\end{aligned}
$$

The index of the root lattice in the weight lattice. $\quad|X: Q|=1$.
The symmetrising matrix $D=\operatorname{diag}\left(d_{i}\right)$.

$$
d_{1}=1, \quad d_{2}=3
$$

The standard invariant form on $H_{\mathbb{R}}$.

$$
\left\langle h_{i}, h_{j}\right\rangle=A_{i j} d_{j} .
$$

The standard invariant form on $H_{\mathbb{R}}{ }^{*}$.

$$
\left\langle\alpha_{i}, \alpha_{j}\right\rangle=d_{i}^{-1} A_{i j}
$$

The Killing form on $H_{\mathbb{R}}$.

$$
\langle x, y\rangle_{K}=\frac{1}{b}\langle x, y\rangle
$$

where $b=8$.

## DYNKIN NAME $\tilde{A}_{1}$ KAC NAME $\tilde{A}_{1}$

Dynkin diagram with labelling.


Generalised Cartan matrix.

$$
\begin{gathered}
0 \\
0\left(\begin{array}{rr}
1 \\
1 & -2 \\
-2 & 2
\end{array}\right)
\end{gathered}
$$

The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+h_{1} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+\alpha_{1} .
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=h_{1} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=\alpha_{1} .
$$

The Coxeter number. $\quad h=2$.
The dual Coxeter number. $\quad h^{v}=2$.
The Lie algebra $L^{0} . \quad L^{0}=A_{1}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1} .
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
A=\left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{1}(h)>0, \alpha_{1}(h)<1\right\}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
A^{*}=\left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{1}\right)>0, \lambda\left(h_{1}\right)<1\right\}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, k \neq 0\} \quad \text { Multiplicity } 1
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0,1$ in terms of the fundamental weights $\bar{\omega}_{i}, \quad i=1$, of $L^{0}$.

$$
\omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+\gamma
$$

The standard invariant form on $H$.

$$
\begin{aligned}
\left\langle h_{i}, h_{j}\right\rangle & =A_{i j} \quad i, j=0,1 \\
\left\langle h_{0}, d\right\rangle & =1, \quad\left\langle h_{1}, d\right\rangle=0 \\
\langle d, d\rangle & =0
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
\left\langle\alpha_{i}, \alpha_{j}\right\rangle & =A_{i j} \quad i, j=0,1 \\
\left\langle\alpha_{0}, \gamma\right\rangle & =1, \quad\left\langle\alpha_{1}, \gamma\right\rangle=0 \\
\langle\gamma, \gamma\rangle & =0
\end{aligned}
$$

## DYNKIN NAME $\tilde{A}_{1}^{\prime} \quad$ KAC NAME ${ }^{2} \tilde{A}_{2}$

 (1st description)Dynkin diagram with labelling.


Generalised Cartan matrix.

$$
\begin{gathered}
0 \\
0\left(\begin{array}{cc}
1 \\
1 \\
2 \\
-4 & -1 \\
2
\end{array}\right)
\end{gathered}
$$

The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=2 h_{0}+h_{1} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+2 \alpha_{1}
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=\frac{1}{2} h_{1} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$

$$
\theta=2 \alpha_{1}
$$

The Coxeter number. $\quad h=3$.
The dual Coxeter number. $\quad h^{v}=3$.
The Lie algebra $L^{0} . \quad L^{0}=A_{1}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\frac{1}{2} \mathbb{Z} h_{1}
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$

$$
M^{*}=\mathbb{Z} \alpha_{1}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
A=\left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{1}(h)>0,2 \alpha_{1}(h)<1\right\}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$

$$
A^{*}=\left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{1}\right)>0, \lambda\left(h_{1}\right)<1\right\}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}, \mathrm{~s}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, \quad r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, \mathrm{l}}=\left\{2 \alpha+(2 r+1) \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, \quad k \neq 0\} \quad \text { Multiplicity } 1
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0,1, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i} \quad i=1, \ldots, l$ of $L^{0}$.

$$
\omega_{0}=2 \gamma, \quad \omega_{1}=\bar{\omega}_{1}+\gamma
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0,1 \\
& \left\langle h_{0}, d\right\rangle=\frac{1}{2}, \quad\left\langle h_{1}, d\right\rangle=0 \\
& \langle d, d\rangle=0
\end{aligned}
$$

The standard invariant form on $H^{*}$

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0,1 \\
& \left\langle\alpha_{0}, \gamma\right\rangle=2, \quad\left\langle\alpha_{1}, \gamma\right\rangle=0 \\
& \langle\gamma, \gamma\rangle=0
\end{aligned}
$$

## DYNKIN NAME $\tilde{A}_{1}^{\prime} \quad$ KAC NAME $\quad{ }^{2} \tilde{A}_{2}$ (2nd description)

Dynkin diagram with labelling.


Generalised Cartan matrix.

$$
\left.\begin{array}{c} 
\\
0 \\
1
\end{array} \begin{array}{rr}
0 & 1 \\
2 & -4 \\
-1 & 2
\end{array}\right)
$$

The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+2 h_{1} .
$$

The basic imaginary root $\delta$.

$$
\delta=2 \alpha_{0}+\alpha_{1}
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=h_{1} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=\alpha_{1}
$$

The Coxeter number. $\quad h=3$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=3$.
The Lie algebra $L^{0} . \quad L^{0}=A_{1}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\frac{1}{2} \mathbb{Z} \alpha_{1}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
A=\left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{1}(h)>0, \alpha_{1}(h)<1\right\}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
A^{*}=\left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{1}\right)>0,2 \lambda\left(h_{1}\right)<1\right\}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}, \mathrm{~s}}=\left\{\frac{1}{2}(\alpha+(2 r-1) \delta) ; \quad \alpha \in \Phi^{0}, \quad r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, \mathrm{l}}=\left\{\alpha+2 r \delta ; \quad \alpha \in \Phi^{0}, \quad r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, \quad k \neq 0\} \text { Multiplicity } 1
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0,1, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i} \quad i=1, \ldots, l$ of $L^{0}$.

$$
\omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+2 \gamma
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0,1 \\
& \left\langle h_{0}, d\right\rangle=2, \quad\left\langle h_{1}, d\right\rangle=0 \\
& \langle d, d\rangle=0
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0,1 \\
& \left\langle\alpha_{0}, \gamma\right\rangle=\frac{1}{2}, \quad\left\langle\alpha_{1}, \gamma\right\rangle=0 \\
& \langle\gamma, \gamma\rangle=0
\end{aligned}
$$

## DYNKIN NAME $\begin{gathered}\tilde{A}_{l} \\ \text { KAC NAME } \\ \tilde{A}_{l} \quad l \geq 2\end{gathered}$

Dynkin diagram with labelling.


Generalised Cartan matrix.


The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+h_{1}+\cdots+h_{l-1}+h_{l}
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+\alpha_{1}+\cdots+\alpha_{l-1}+\alpha_{l} .
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=h_{1}+h_{2}+\cdots+h_{l-1}+h_{l} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=\alpha_{1}+\alpha_{2}+\cdots+\alpha_{l-1}+\alpha_{l}
$$

The Coxeter number. $\quad h=l+1$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=l+1$.
The Lie algebra $L^{0} . \quad L^{0}=A_{l}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+\cdots+\mathbb{Z} h_{l-1}+\mathbb{Z} h_{l}
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+\cdots+\mathbb{Z} \alpha_{l-1}+\mathbb{Z} \alpha_{l}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A=\{ & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, l\right. \\
& \left.\alpha_{1}(h)+\alpha_{2}(h)+\cdots+\alpha_{l-1}(h)+\alpha_{l}(h)<1\right\}
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}= & \left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, l\right. \\
& \left.\lambda\left(h_{1}\right)+\lambda\left(h_{2}\right)+\cdots+\lambda\left(h_{l-1}\right)+\lambda\left(h_{l}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, k \neq 0\} \text { Multiplicity } l .
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0,1, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i}, \quad i=1, \ldots, l$ of $L^{0}$.
$\omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+\gamma, \quad \omega_{2}=\bar{\omega}_{2}+\gamma, \quad \ldots, \quad \omega_{l-1}=\bar{\omega}_{l-1}+\gamma, \quad \omega_{l}=\bar{\omega}_{l}+\gamma$.

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, l \\
& \langle d, d\rangle=0 .
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, l \\
& \langle\gamma, \gamma\rangle=0
\end{aligned}
$$

DYNKIN NAME $\begin{gathered}\tilde{B}_{l} \\ \text { KAC NAME }\end{gathered} \tilde{B}_{l} \quad l \geq 3$
Dynkin diagram with labelling.


Generalised Cartan matrix.


The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+h_{1}+2 h_{2}+2 h_{3}+\cdots+2 h_{l-1}+h_{l} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+\alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\cdots+2 \alpha_{l-1}+2 \alpha_{l} .
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=h_{1}+2 h_{2}+2 h_{3}+\cdots+2 h_{l-1}+h_{l} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=\alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\cdots+2 \alpha_{l-1}+2 \alpha_{l} .
$$

The Coxeter number. $\quad h=2 l$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=2 l-1$.
The Lie algebra $L^{0} . \quad L^{0}=B_{l}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+\cdots+\mathbb{Z} h_{l-1}+\mathbb{Z} h_{l}
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+\cdots+\mathbb{Z} \alpha_{l-1}+2 \mathbb{Z} \alpha_{l}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, l\right. \\
& \left.\alpha_{1}(h)+2 \alpha_{2}(h)+2 \alpha_{3}(h)+\cdots+2 \alpha_{l-1}(h)+2 \alpha_{l}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}= & \left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, l\right. \\
& \left.\lambda\left(h_{1}\right)+2 \lambda\left(h_{2}\right)+2 \lambda\left(h_{3}\right)+\cdots+2 \lambda\left(h_{l-1}\right)+\lambda\left(h_{l}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, k \neq 0\} \quad \text { Multiplicity } l .
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i}, \quad i=1, \ldots, l$ of $L^{0}$.

$$
\omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+\gamma, \quad \omega_{2}=\bar{\omega}_{2}+2 \gamma, \ldots, \omega_{l-1}=\bar{\omega}_{l-1}+2 \gamma, \quad \omega_{l}=\bar{\omega}_{l}+\gamma .
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, l \\
& \langle d, d\rangle=0
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, l \\
& \langle\gamma, \gamma\rangle=0
\end{aligned}
$$

## DYNKIN NAME $\tilde{B}_{l}^{\mathrm{t}} \quad$ KAC NAME ${ }^{2} \tilde{A}_{2 l-1} \quad l \geq 3$

Dynkin diagram with labelling.


Generalised Cartan matrix.


The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+h_{1}+2 h_{2}+2 h_{3}+\cdots+2 h_{l-1}+2 h_{l} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+\alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\cdots+2 \alpha_{l-1}+\alpha_{l}
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=h_{1}+2 h_{2}+2 h_{3}+\cdots+2 h_{l-1}+2 h_{l} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=\alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\cdots+2 \alpha_{l-1}+\alpha_{l}
$$

The Coxeter number. $\quad h=2 l-1$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=2 l$.
The Lie algebra $L^{0} . \quad L^{0}=C_{l}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+\cdots+\mathbb{Z} h_{l-1}+2 \mathbb{Z} h_{l}
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+\cdots+\mathbb{Z} \alpha_{l-1}+\mathbb{Z} \alpha_{l}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A=\{ & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, l\right. \\
& \left.\alpha_{1}(h)+2 \alpha_{2}(h)+2 \alpha_{3}(h)+\cdots+2 \alpha_{l-1}(h)+\alpha_{l}(h)<1\right\}
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}= & \left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, l\right. \\
& \left.\lambda\left(h_{1}\right)+2 \lambda\left(h_{2}\right)+2 \lambda\left(h_{3}\right)+\cdots+2 \lambda\left(h_{l-1}\right)+2 \lambda\left(h_{l}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
\Phi_{\mathrm{Re}, \mathrm{~s}}= & \left\{\alpha+r \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}, r \in \mathbb{Z}\right\} \\
\Phi_{\mathrm{Re}, \mathrm{l}}= & \left\{\alpha+2 r \delta ; \alpha \in \Phi_{1}^{0}, r \in \mathbb{Z}\right\} \\
\Phi_{\mathrm{Im}}= & \{2 k \delta ; k \in \mathbb{Z}, k \neq 0\} \quad \text { Multiplicity } l \\
& \cup\{(2 k+1) \delta ; k \in \mathbb{Z}\} \quad \text { Multiplicity } l-1 .
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0,1, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i}, \quad i=1, \ldots, l$ of $L^{0}$.

$$
\begin{aligned}
& \omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+\gamma, \quad \omega_{2}=\bar{\omega}_{2}+2 \gamma, \quad \ldots \\
& \omega_{l-1}=\bar{\omega}_{l-1}+2 \gamma, \quad \omega_{l}=\bar{\omega}_{l}+2 \gamma
\end{aligned}
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, l \\
& \langle d, d\rangle=0 .
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, l \\
& \langle\gamma, \gamma\rangle=0 .
\end{aligned}
$$

## DYNKIN NAME $\begin{gathered}\tilde{C}_{l} \\ \text { KAC NAME }\end{gathered} \tilde{C}_{l} \quad l \geq 2$

Dynkin diagram with labelling.


Generalised Cartan matrix


The integers $a_{0}, a_{1}, \ldots, a_{l}$


The integers $c_{0}, c_{1}, \ldots, c_{l}$


The central element $c$.

$$
c=h_{0}+h_{1}+h_{2}+\cdots+h_{l-1}+h_{l} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+2 \alpha_{1}+2 \alpha_{2}+\cdots+2 \alpha_{l-1}+\alpha_{l}
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=h_{1}+h_{2}+\cdots+h_{l-1}+h_{l} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=2 \alpha_{1}+2 \alpha_{2}+\cdots+2 \alpha_{l-1}+\alpha_{l} .
$$

The Coxeter number. $\quad h=2 l$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=l+1$.
The Lie algebra $L^{0} . \quad L^{0}=C_{l}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+\cdots+\mathbb{Z} h_{l-1}+\mathbb{Z} h_{l} .
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=2 \mathbb{Z} \alpha_{1}+2 \mathbb{Z} \alpha_{2}+\cdots+2 \mathbb{Z} \alpha_{l-1}+\mathbb{Z} \alpha_{l}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, l\right. \\
& \left.2 \alpha_{1}(h)+2 \alpha_{2}(h)+\cdots+2 \alpha_{l-1}(h)+\alpha_{l}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}= & \left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, l\right. \\
& \left.\lambda\left(h_{1}\right)+\lambda\left(h_{2}\right)+\cdots+\lambda\left(h_{l-1}\right)+\lambda\left(h_{l}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, k \neq 0\} \quad \text { Multiplicity } l
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0,1, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i} \quad i=1, \ldots, l$ of $L^{0}$.
$\omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+\gamma, \quad \omega_{2}=\bar{\omega}_{2}+\gamma, \quad \ldots, \quad \omega_{l-1}=\bar{\omega}_{l-1}+\gamma, \quad \omega_{l}=\bar{\omega}_{l}+\gamma$.
The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, l \\
& \langle d, d\rangle=0
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, l \\
& \langle\gamma, \gamma\rangle=0
\end{aligned}
$$

## DYNKIN NAME $\quad \tilde{C}_{l}^{\mathrm{t}} \quad$ KAC NAME $\quad{ }^{2} \tilde{D}_{l+1} \quad l \geq 2$

Dynkin diagram with labelling.


Generalised Cartan matrix.


The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+2 h_{1}+2 h_{2}+\cdots+2 h_{l-1}+h_{l} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+\alpha_{1}+\alpha_{2}+\cdots+\alpha_{l-1}+\alpha_{l} .
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=2 h_{1}+2 h_{2}+\cdots+2 h_{l-1}+h_{l} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=\alpha_{1}+\alpha_{2}+\cdots+\alpha_{l-1}+\alpha_{l}
$$

The Coxeter number. $\quad h=l+1$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=2 l$.
The Lie algebra $L^{0} . \quad L^{0}=B_{l}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=2 \mathbb{Z} h_{1}+2 \mathbb{Z} h_{2}+\cdots+2 \mathbb{Z} h_{l-1}+\mathbb{Z} h_{l}
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+\cdots+\mathbb{Z} \alpha_{l-1}+\mathbb{Z} \alpha_{l}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, l\right. \\
& \left.\alpha_{1}(h)+\alpha_{2}(h)+\cdots+\alpha_{l-1}(h)+\alpha_{l}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}= & \left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, l\right. \\
& \left.2 \lambda\left(h_{1}\right)+2 \lambda\left(h_{2}\right)+\cdots+2 \lambda\left(h_{l-1}\right)+\lambda\left(h_{l}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
\Phi_{\mathrm{Re}, \mathrm{~s}}= & \left\{\alpha+r \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}, \quad r \in \mathbb{Z}\right\} \\
\Phi_{\mathrm{Re}, \mathrm{l}}= & \left\{\alpha+2 r \delta ; \alpha \in \Phi_{1}^{0}, \quad r \in \mathbb{Z}\right\} \\
\Phi_{\mathrm{Im}}= & \{2 k \delta ; k \in \mathbb{Z}, \quad k \neq 0\} \quad \text { Multiplicity } l \\
& \cup\{(2 k+1) \delta ; k \in \mathbb{Z}\} \quad \text { Multiplicity } 1 .
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0,1, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i}, \quad i=1, \ldots, l$ of $L^{0}$.
$\omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+2 \gamma, \quad \omega_{2}=\bar{\omega}_{2}+2 \gamma, \quad \ldots, \omega_{l-1}=\bar{\omega}_{l-1}+2 \gamma, \quad \omega_{l}=\bar{\omega}_{l}+\gamma$.

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, l \\
& \langle d, d\rangle=0 .
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, l \\
& \langle\gamma, \gamma\rangle=0 .
\end{aligned}
$$

## DYNKIN NAME $\quad \tilde{C}_{l}^{\prime} \quad$ KAC NAME $\quad{ }^{2} \tilde{A}_{2 l} \quad l \geq 2$ <br> (1st description)

Dynkin diagram with labelling.


Generalised Cartan matrix.


The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+2 h_{1}+2 h_{2}+\cdots+2 h_{l-1}+2 h_{l}
$$

The basic imaginary root $\delta$.

$$
\delta=2 \alpha_{0}+2 \alpha_{1}+2 \alpha_{2}+\cdots+2 \alpha_{l-1}+\alpha_{l}
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=h_{1}+h_{2}+\cdots+h_{l-1}+h_{l} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=2 \alpha_{1}+2 \alpha_{2}+\cdots+2 \alpha_{l-1}+\alpha_{l}
$$

The Coxeter number. $\quad h=2 l+1$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=2 l+1$.
The Lie algebra $L^{0} . \quad L^{0}=C_{l}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+\cdots+\mathbb{Z} h_{l-1}+\mathbb{Z} h_{l} .
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+\cdots+\mathbb{Z} \alpha_{l-1}+\frac{1}{2} \mathbb{Z} \alpha_{l}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, l\right. \\
& \left.2 \alpha_{1}(h)+2 \alpha_{2}(h)+\cdots+2 \alpha_{l-1}(h)+\alpha_{l}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}= & \left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, l\right. \\
& \left.2 \lambda\left(h_{1}\right)+2 \lambda\left(h_{2}\right)+\cdots+2 \lambda\left(h_{l-1}\right)+2 \lambda\left(h_{l}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}, \mathrm{~s}}=\left\{\frac{1}{2}(\alpha+(2 r-1) \delta) ; \quad \alpha \in \Phi_{l}^{0}, \quad r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, \mathrm{i}}=\left\{\alpha+r \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, \mathrm{l}}=\left\{\alpha+2 r \delta ; \alpha \in \Phi_{1}^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, \quad k \neq 0\} \quad \text { Multiplicity } l .
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0,1, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i} \quad i=1, \ldots, l$ of $L^{0}$.

$$
\begin{aligned}
& \omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+2 \gamma, \quad \omega_{2}=\bar{\omega}_{2}+2 \gamma, \quad \ldots, \\
& \omega_{l-1}=\bar{\omega}_{l-1}+2 \gamma, \quad \omega_{l}=\bar{\omega}_{l}+2 \gamma .
\end{aligned}
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle h_{0}, d\right\rangle=2, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, l \\
& \langle d, d\rangle=0
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=\frac{1}{2}, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, l \\
& \langle\gamma, \gamma\rangle=0
\end{aligned}
$$

DYNKIN NAME $\quad \tilde{C}_{l}^{\prime} \quad$ KAC NAME $\quad{ }^{2} \tilde{A}_{2 l} \quad l \geq 2$
(2nd description)
Dynkin diagram with labelling.


Generalised Cartan matrix.


The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=2 h_{0}+2 h_{1}+2 h_{2}+\cdots+2 h_{l-1}+h_{l} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+2 \alpha_{1}+2 \alpha_{2}+\cdots+2 \alpha_{l-1}+2 \alpha_{l} .
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=h_{1}+h_{2}+\cdots+h_{l-1}+\frac{1}{2} h_{l}
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=2 \alpha_{1}+2 \alpha_{2}+\cdots+2 \alpha_{l-1}+2 \alpha_{l}
$$

The Coxeter number. $\quad h=2 l+1$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=2 l+1$.
The Lie algebra $L^{0} . \quad L^{0}=B_{l}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+\cdots+\mathbb{Z} h_{l-1}+\frac{1}{2} \mathbb{Z} h_{l} .
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+\cdots+\mathbb{Z} \alpha_{l-1}+\mathbb{Z} \alpha_{l}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, l\right. \\
& \left.2 \alpha_{1}(h)+2 \alpha_{2}(h)+\cdots+2 \alpha_{l-1}(h)+2 \alpha_{l}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}= & \left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, l\right. \\
& \left.2 \lambda\left(h_{1}\right)+2 \lambda\left(h_{2}\right)+\cdots+2 \lambda\left(h_{l-1}\right)+\lambda\left(h_{l}\right)<1\right\}
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}, \mathrm{~s}}=\left\{\alpha+r \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}, \quad r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, \mathrm{i}}=\left\{\alpha+r \delta ; \alpha \in \Phi_{1}^{0}, \quad r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Re}, \mathrm{l}}=\left\{2 \alpha+(2 r+1) \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}, \quad r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, \quad k \neq 0\} \quad \text { Multiplicity } l .
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0,1, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i} \quad i=1, \ldots, l$ of $L^{0}$.

$$
\begin{aligned}
& \omega_{0}=2 \gamma, \quad \omega_{1}=\bar{\omega}_{1}+2 \gamma, \quad \omega_{2}=\bar{\omega}_{2}+2 \gamma, \quad \ldots \\
& \omega_{l-1}=\bar{\omega}_{l-1}+2 \gamma, \quad \omega_{l}=\bar{\omega}_{l}+\gamma
\end{aligned}
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle h_{0}, d\right\rangle=\frac{1}{2}, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, l \\
& \langle d, d\rangle=0
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=2, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, l \\
& \langle\gamma, \gamma\rangle=0
\end{aligned}
$$

## $\begin{array}{lllll}\text { DYNKIN NAME } & \tilde{D}_{4} & \text { KAC NAME } & \tilde{D}_{4}\end{array}$

Dynkin diagram with labelling.


Generalised Cartan matrix.

$$
\begin{aligned}
& 0 \\
& 1 \\
& 2 \\
& 3 \\
& 4
\end{aligned}\left(\begin{array}{rrrrr}
0 & 1 & 2 & 3 & 4 \\
2 & 0 & -1 & 0 & 0 \\
0 & 2 & -1 & 0 & 0 \\
-1 & -1 & 2 & -1 & -1 \\
0 & 0 & -1 & 2 & 0 \\
0 & 0 & -1 & 0 & 2
\end{array}\right)
$$

The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+h_{1}+2 h_{2}+h_{3}+h_{4} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+\alpha_{1}+2 \alpha_{2}+\alpha_{3}+\alpha_{4} .
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=h_{1}+2 h_{2}+h_{3}+h_{4} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=\alpha_{1}+2 \alpha_{2}+\alpha_{3}+\alpha_{4}
$$

The Coxeter number. $\quad h=6$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=6$.

The Lie algebra $L^{0} . \quad L^{0}=D_{4}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+\mathbb{Z} h_{3}+\mathbb{Z} h_{4} .
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+\mathbb{Z} \alpha_{3}+\mathbb{Z} \alpha_{4}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, 4\right. \\
& \left.\alpha_{1}(h)+2 \alpha_{2}(h)+\alpha_{3}(h)+\alpha_{4}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}= & \left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, l\right. \\
& \left.\lambda\left(h_{1}\right)+2 \lambda\left(h_{2}\right)+\lambda\left(h_{3}\right)+\lambda\left(h_{4}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, k \neq 0\} \quad \text { Multiplicity } 4 .
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0,1, \ldots l$ in terms of the fundamental weights $\overline{\omega_{i}} \quad i=1, \ldots, l$ of $L^{0}$.

$$
\omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+\gamma, \quad \omega_{2}=\bar{\omega}_{2}+2 \gamma, \quad \omega_{3}=\bar{\omega}_{3}+\gamma, \quad \omega_{4}=\bar{\omega}_{4}+\gamma .
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=A_{i j} \quad i, j=0,1,2,3,4 \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1,2,3,4 \\
& \langle d, d\rangle=0 .
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=A_{i j} \quad i, j=0,1,2,3,4 \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1,2,3,4 \\
& \langle\gamma, \gamma\rangle=0 .
\end{aligned}
$$

## DYNKIN NAME $\tilde{D}_{l} \quad$ KAC NAME $\quad \tilde{D}_{l} \quad l \geq 5$

Dynkin diagram with labelling.


Generalised Cartan matrix.


The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+h_{1}+2 h_{2}+2 h_{3}+\cdots+2 h_{l-2}+h_{l-1}+h_{l} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+\alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\cdots+2 \alpha_{l-2}+\alpha_{l-1}+\alpha_{l}
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=h_{1}+2 h_{2}+2 h_{3}+\cdots+2 h_{l-2}+h_{l-1}+h_{l} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=\alpha_{1}+2 \alpha_{2}+2 \alpha_{3}+\cdots+2 \alpha_{l-2}+\alpha_{l-1}+\alpha_{l}
$$

The Coxeter number. $\quad h=2 l-2$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=2 l-2$.
The Lie algebra $L^{0} . \quad L^{0}=D_{l}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+\cdots+\mathbb{Z} h_{l-1}+\mathbb{Z} h_{l}
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+\cdots+\mathbb{Z} \alpha_{l-1}+\mathbb{Z} \alpha_{l}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, l\right. \\
& \left.\alpha_{1}(h)+2 \alpha_{2}(h)+2 \alpha_{3}(h)+\cdots+2 \alpha_{l-2}(h)+\alpha_{l-1}(h)+\alpha_{l}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}= & \left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, l\right. \\
& \left.\lambda\left(h_{1}\right)+2 \lambda\left(h_{2}\right)+2 \lambda\left(h_{3}\right)+\cdots+2 \lambda\left(h_{l-2}\right)+\lambda\left(h_{l-1}\right)+\lambda\left(h_{l}\right)<1\right\}
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, \quad k \neq 0\} \quad \text { Multiplicity } l
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*}, \quad i=0,1, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i}, \quad i=1, \ldots, l$ of $L^{0}$.

$$
\begin{aligned}
& \omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+\gamma, \quad \omega_{2}=\bar{\omega}_{2}+2 \gamma, \quad \ldots \\
& \omega_{l-2}=\bar{\omega}_{l-2}+2 \gamma, \quad \omega_{l-1}=\bar{\omega}_{l-1}+\gamma, \quad \omega_{l}=\bar{\omega}_{l}+\gamma
\end{aligned}
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, l \\
& \langle d, d\rangle=0
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=A_{i j} \quad i, j=0,1, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, l \\
& \langle\gamma, \gamma\rangle=0
\end{aligned}
$$

## DYNKIN NAME $\quad \tilde{E}_{6} \quad$ KAC NAME $\quad \tilde{E}_{6}$

Dynkin diagram with labelling.


Generalised Cartan matrix.
0
1
2
3
4
5
6 $\left(\begin{array}{rrrrrrr}0 & 1 & 2 & 3 & 4 & 5 & 6 \\ 2 & 0 & 0 & 0 & -1 & 0 & 0 \\ 0 & 2 & -1 & 0 & 0 & 0 & 0 \\ 0 & -1 & 2 & -1 & 0 & 0 & 0 \\ 0 & 0 & -1 & 2 & -1 & -1 & 0 \\ -1 & 0 & 0 & -1 & 2 & 0 & 0 \\ 0 & 0 & 0 & -1 & 0 & 2 & -1 \\ 0 & 0 & 0 & 0 & 0 & -1 & 2\end{array}\right)$

The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+h_{1}+2 h_{2}+3 h_{3}+2 h_{4}+2 h_{5}+h_{6} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+\alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+2 \alpha_{4}+2 \alpha_{5}+\alpha_{6}
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=h_{1}+2 h_{2}+3 h_{3}+2 h_{4}+2 h_{5}+h_{6} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=\alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+2 \alpha_{4}+2 \alpha_{5}+\alpha_{6}
$$

The Coxeter number. $\quad h=12$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=12$.
The Lie algebra $L^{0} . \quad L^{0}=E_{6}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+\mathbb{Z} h_{3}+\mathbb{Z} h_{4}+\mathbb{Z} h_{5}+\mathbb{Z} h_{6}
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+\mathbb{Z} \alpha_{3}+\mathbb{Z} \alpha_{4}+\mathbb{Z} \alpha_{5}+\mathbb{Z} \alpha_{6}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, 6\right. \\
& \left.\alpha_{1}(h)+2 \alpha_{2}(h)+3 \alpha_{3}(h)+2 \alpha_{4}(h)+2 \alpha_{5}(h)+\alpha_{6}(h)<1\right\}
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}= & \left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, 6\right. \\
& \left.\lambda\left(h_{1}\right)+2 \lambda\left(h_{2}\right)+3 \lambda\left(h_{3}\right)+2 \lambda\left(h_{4}\right)+2 \lambda\left(h_{5}\right)+\lambda\left(h_{6}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, k \neq 0\} \quad \text { Multiplicity } 6 .
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*}, \quad i=0, \ldots, 6$ in terms of the fundamental weights $\bar{\omega}_{i}, \quad i=1, \ldots, 6$ of $L^{0}$.

$$
\begin{aligned}
& \omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+\gamma, \quad \omega_{2}=\bar{\omega}_{2}+2 \gamma, \quad \omega_{3}=\bar{\omega}_{3}+3 \gamma \\
& \omega_{4}=\bar{\omega}_{4}+2 \gamma, \quad \omega_{5}=\bar{\omega}_{5}+2 \gamma, \quad \omega_{6}=\bar{\omega}_{6}+\gamma
\end{aligned}
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=A_{i j} \quad i, j=0, \ldots, 6 \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, 6 \\
& \langle d, d\rangle=0 .
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=A_{i j} \quad i, j=0, \ldots, 6 \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, 6 \\
& \langle\gamma, \gamma\rangle=0 .
\end{aligned}
$$

## DYNKIN NAME $\quad \tilde{E}_{7} \quad$ KAC NAME $\quad \tilde{E}_{7}$

Dynkin diagram with labelling.


Generalised Cartan matrix.
0
1
2
3
4
5
6
7 $\left(\begin{array}{rrrrrrrr}0 & 1 & 2 & 3 & 4 & 5 & 6 & 7 \\ 2 & 0 & 0 & 0 & 0 & 0 & 0 & -1 \\ 0 & 2 & -1 & 0 & 0 & 0 & 0 & 0 \\ 0 & -1 & 2 & -1 & 0 & 0 & 0 & 0 \\ 0 & 0 & -1 & 2 & -1 & 0 & 0 & 0 \\ 0 & 0 & 0 & -1 & 2 & -1 & -1 & 0 \\ 0 & 0 & 0 & 0 & -1 & 2 & 0 & 0 \\ 0 & 0 & 0 & 0 & -1 & 0 & 2 & -1 \\ -1 & 0 & 0 & 0 & 0 & 0 & -1 & 2\end{array}\right)$

The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+h_{1}+2 h_{2}+3 h_{3}+4 h_{4}+2 h_{5}+3 h_{6}+2 h_{7} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+\alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+4 \alpha_{4}+2 \alpha_{5}+3 \alpha_{6}+2 \alpha_{7}
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=h_{1}+2 h_{2}+3 h_{3}+4 h_{4}+2 h_{5}+3 h_{6}+2 h_{7} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=\alpha_{1}+2 \alpha_{2}+3 \alpha_{3}+4 \alpha_{4}+2 \alpha_{5}+3 \alpha_{6}+2 \alpha_{7} .
$$

The Coxeter number. $\quad h=18$.
The dual Coxeter number. $\quad h^{v}=18$.
The Lie algebra $L^{0} . \quad L^{0}=E_{7}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+\mathbb{Z} h_{3}+\mathbb{Z} h_{4}+\mathbb{Z} h_{5}+\mathbb{Z} h_{6}+\mathbb{Z} h_{7} .
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+\mathbb{Z} \alpha_{3}+\mathbb{Z} \alpha_{4}+\mathbb{Z} \alpha_{5}+\mathbb{Z} \alpha_{6}+\mathbb{Z} \alpha_{7} .
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, 7\right. \\
& \left.\alpha_{1}(h)+2 \alpha_{2}(h)+3 \alpha_{3}(h)+4 \alpha_{4}(h)+2 \alpha_{5}(h)+3 \alpha_{6}(h)+2 \alpha_{7}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}=\{ & \lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, 7 \\
& \left.\lambda\left(h_{1}\right)+2 \lambda\left(h_{2}\right)+3 \lambda\left(h_{3}\right)+4 \lambda\left(h_{4}\right)+2 \lambda\left(h_{5}\right)+3 \lambda\left(h_{6}\right)+2 \lambda\left(h_{7}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, k \neq 0\} \quad \text { Multiplicity } 7 .
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*}, \quad i=0, \ldots, 7$ in terms of the fundamental weights $\bar{\omega}_{i}, \quad i=1, \ldots, 7$ of $L^{0}$.

$$
\begin{array}{lll}
\omega_{0}=\gamma, & \omega_{1}=\bar{\omega}_{1}+\gamma, & \omega_{2}=\bar{\omega}_{2}+2 \gamma, \\
\omega_{5}=\bar{\omega}_{5}+2 \gamma, & \omega_{6}=\bar{\omega}_{3}+3 \gamma, & \omega_{4}=\bar{\omega}_{4}+4 \gamma, \\
\omega_{7}=\bar{\omega}_{7}+2 \gamma .
\end{array}
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=A_{i j} \quad i, j=0, \ldots, 7 \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, 7 \\
& \langle d, d\rangle=0
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=A_{i j} \quad i, j=0, \ldots, 7 \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, 7 \\
& \langle\gamma, \gamma\rangle=0
\end{aligned}
$$

## DYNKIN NAME $\tilde{E}_{8}$ KAC NAME $\tilde{E}_{8}$

Dynkin diagram with labelling.


Generalised Cartan matrix.
0
1
2
3
4
5
6
7
8 $\left(\begin{array}{rrrrrrrrr}0 & 1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 \\ 2 & -1 & 0 & 0 & 0 & 0 & 0 & 0 & 0 \\ -1 & 2 & -1 & 0 & 0 & 0 & 0 & 0 & 0 \\ 0 & -1 & 2 & -1 & 0 & 0 & 0 & 0 & 0 \\ 0 & 0 & -1 & 2 & -1 & 0 & 0 & 0 & 0 \\ 0 & 0 & 0 & -1 & 2 & -1 & 0 & 0 & 0 \\ 0 & 0 & 0 & 0 & -1 & 2 & -1 & -1 & 0 \\ 0 & 0 & 0 & 0 & 0 & -1 & 2 & 0 & 0 \\ 0 & 0 & 0 & 0 & 0 & -1 & 0 & 2 & -1 \\ 0 & 0 & 0 & 0 & 0 & 0 & 0 & -1 & 2\end{array}\right)$

The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+2 h_{1}+3 h_{2}+4 h_{3}+5 h_{4}+6 h_{5}+3 h_{6}+4 h_{7}+2 h_{8} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+2 \alpha_{1}+3 \alpha_{2}+4 \alpha_{3}+5 \alpha_{4}+6 \alpha_{5}+3 \alpha_{6}+4 \alpha_{7}+2 \alpha_{8}
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=2 h_{1}+3 h_{2}+4 h_{3}+5 h_{4}+6 h_{5}+3 h_{6}+4 h_{7}+2 h_{8}
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=2 \alpha_{1}+3 \alpha_{2}+4 \alpha_{3}+5 \alpha_{4}+6 \alpha_{5}+3 \alpha_{6}+4 \alpha_{7}+2 \alpha_{8}
$$

The Coxeter number. $\quad h=30$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=30$.
The Lie algebra $L^{0} . \quad L^{0}=E_{8}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+\mathbb{Z} h_{3}+\mathbb{Z} h_{4}+\mathbb{Z} h_{5}+\mathbb{Z} h_{6}+\mathbb{Z} h_{7}+\mathbb{Z} h_{8}
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+\mathbb{Z} \alpha_{3}+\mathbb{Z} \alpha_{4}+\mathbb{Z} \alpha_{5}+\mathbb{Z} \alpha_{6}+\mathbb{Z} \alpha_{7}+\mathbb{Z} \alpha_{8}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A=\{ & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, 8\right. \\
& 2 \alpha_{1}(h)+3 \alpha_{2}(h)+4 \alpha_{3}(h)+5 \alpha_{4}(h)+6 \alpha_{5}(h) \\
& \left.+3 \alpha_{6}(h)+4 \alpha_{7}(h)+2 \alpha_{8}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}=\{ & \lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, 8 \\
& 2 \lambda\left(h_{1}\right)+3 \lambda\left(h_{2}\right)+4 \lambda\left(h_{3}\right)+5 \lambda\left(h_{4}\right)+6 \lambda\left(h_{5}\right) \\
& \left.+3 \lambda\left(h_{6}\right)+4 \lambda\left(h_{7}\right)+2 \lambda\left(h_{8}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, k \neq 0\} \quad \text { Multiplicity } 8
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*}, \quad i=0, \ldots, 8$ in terms of the fundamental weights $\bar{\omega}_{i}, \quad i=1, \ldots, 8$ of $L^{0}$.

$$
\begin{aligned}
& \omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+2 \gamma, \quad \omega_{2}=\bar{\omega}_{2}+3 \gamma, \quad \omega_{3}=\bar{\omega}_{3}+4 \gamma, \quad \omega_{4}=\bar{\omega}_{4}+5 \gamma \\
& \omega_{5}=\bar{\omega}_{5}+6 \gamma, \quad \omega_{6}=\bar{\omega}_{6}+3 \gamma, \quad \omega_{7}=\bar{\omega}_{7}+4 \gamma, \quad \omega_{8}=\bar{\omega}_{8}+2 \gamma
\end{aligned}
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=A_{i j} \quad i, j=0, \ldots, 8 \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, 8 \\
& \langle d, d\rangle=0
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=A_{i j} \quad i, j=0, \ldots, 8 \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, 8 \\
& \langle\gamma, \gamma\rangle=0
\end{aligned}
$$

## DYNKIN NAME $\tilde{F}_{4}$ KAC NAME $\tilde{F}_{4}$

Dynkin diagram with labelling.


Generalised Cartan matrix.

$$
\begin{aligned}
& 0 \\
& 1 \\
& 2 \\
& 3 \\
& 4
\end{aligned}\left(\begin{array}{rrrrr}
0 & 1 & 2 & 3 & 4 \\
2 & -1 & 0 & 0 & 0 \\
-1 & 2 & -1 & 0 & 0 \\
0 & -1 & 2 & -1 & 0 \\
0 & 0 & -2 & 2 & -1 \\
0 & 0 & 0 & -1 & 2
\end{array}\right)
$$

The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+2 h_{1}+3 h_{2}+2 h_{3}+h_{4}
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+2 \alpha_{1}+3 \alpha_{2}+4 \alpha_{3}+2 \alpha_{4}
$$

The element $h_{0} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=2 h_{1}+3 h_{2}+2 h_{3}+h_{4}
$$

The element $\theta \in\left(h_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=2 \alpha_{1}+3 \alpha_{2}+4 \alpha_{3}+2 \alpha_{4}
$$

The Coxeter number. $\quad h=12$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=9$.
The Lie algebra $L^{0} . \quad L^{0}=F_{4}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+\mathbb{Z} h_{3}+\mathbb{Z} h_{4}
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+2 \mathbb{Z} \alpha_{3}+2 \mathbb{Z} \alpha_{4} .
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, 4\right. \\
& \left.2 \alpha_{1}(h)+3 \alpha_{2}(h)+4 \alpha_{3}(h)+2 \alpha_{4}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$

$$
\begin{aligned}
A^{*}=\{ & \lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, 4 \\
& \left.2 \lambda\left(h_{1}\right)+3 \lambda\left(h_{2}\right)+2 \lambda\left(h_{3}\right)+\lambda\left(h_{4}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, k \neq 0\} \quad \text { Multiplicity } 4 .
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i}, \quad i=1, \ldots, l$ of $L^{0}$.

$$
\omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+2 \gamma, \quad \omega_{2}=\bar{\omega}_{2}+3 \gamma, \quad \omega_{3}=\bar{\omega}_{3}+2 \gamma, \quad \omega_{4}=\bar{\omega}_{4}+\gamma .
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0, \ldots, l \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, l \\
& \langle d, d\rangle=0 .
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, l \\
& \langle\gamma, \gamma\rangle=0 .
\end{aligned}
$$

## DYNKIN NAME $\begin{array}{llll} & \tilde{F}_{4}^{\mathrm{t}} & \text { KAC NAME } & \\ & { }^{2} \tilde{E}_{6}\end{array}$

Dynkin diagram with labelling.


Generalised Cartan matrix.

$$
\begin{aligned}
& \\
& 0 \\
& 1 \\
& 2 \\
& 3 \\
& 4
\end{aligned}\left(\begin{array}{rrrrr}
0 & 1 & 2 & 3 & 4 \\
2 & -1 & 0 & 0 & 0 \\
-1 & 2 & -1 & 0 & 0 \\
0 & -1 & 2 & -2 & 0 \\
0 & 0 & -1 & 2 & -1 \\
0 & 0 & 0 & -1 & 2
\end{array}\right)
$$

The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+2 h_{1}+3 h_{2}+4 h_{3}+2 h_{4} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+2 \alpha_{1}+3 \alpha_{2}+2 \alpha_{3}+\alpha_{4}
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=2 h_{1}+3 h_{2}+4 h_{3}+2 h_{4} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=2 \alpha_{1}+3 \alpha_{2}+2 \alpha_{3}+\alpha_{4}
$$

The Coxeter number. $\quad h=9$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=12$.
The Lie algebra $L^{0} . \quad L^{0}=F_{4}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2}+2 \mathbb{Z} h_{3}+2 \mathbb{Z} h_{4}
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2}+\mathbb{Z} \alpha_{3}+\mathbb{Z} \alpha_{4} .
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1, \ldots, 4\right. \\
& \left.2 \alpha_{1}(h)+3 \alpha_{2}(h)+2 \alpha_{3}(h)+\alpha_{4}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}= & \left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1, \ldots, 4\right. \\
& \left.2 \lambda\left(h_{1}\right)+3 \lambda\left(h_{2}\right)+4 \lambda\left(h_{3}\right)+2 \lambda\left(h_{4}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
\Phi_{\mathrm{Re}, \mathrm{~s}}= & \left\{\alpha+r \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}, r \in \mathbb{Z}\right\} \\
\Phi_{\mathrm{Re},, 1}= & \left\{\alpha+2 r \delta ; \alpha \in \Phi_{l}^{0}, r \in \mathbb{Z}\right\} \\
\Phi_{\mathrm{Im}}= & \{2 k \delta ; k \in \mathbb{Z}, k \neq 0\} \quad \text { Multiplicity 4 } \\
& \cup\{(2 k+1) \delta ; k \in \mathbb{Z}\} \quad \text { Multiplicity 2. }
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i}, \quad i=1, \ldots, l$ of $L^{0}$.

$$
\omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+2 \gamma, \quad \omega_{2}=\bar{\omega}_{2}+3 \gamma, \quad \omega_{3}=\bar{\omega}_{3}+4 \gamma, \quad \omega_{4}=\bar{\omega}_{4}+2 \gamma .
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0, \ldots l \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots l \\
& \langle d, d\rangle=0 .
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots l \\
& \langle\gamma, \gamma\rangle=0 .
\end{aligned}
$$

## DYNKIN NAME $\quad \tilde{\sigma}_{2} \quad$ KAC NAME $\quad \tilde{G}_{2}$

Dynkin diagram with labelling.


Generalised Cartan matrix.

$$
\begin{aligned}
& \\
& 0 \\
& 1 \\
& 2
\end{aligned}\left(\begin{array}{rrr}
0 & 1 & 2 \\
2 & -1 & 0 \\
-1 & 2 & -1 \\
0 & -3 & 2
\end{array}\right)
$$

The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+2 h_{1}+h_{2} .
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+2 \alpha_{1}+3 \alpha_{2}
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=2 h_{1}+h_{2} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=2 \alpha_{1}+3 \alpha_{2} .
$$

The Coxeter number. $\quad h=6$.
The dual Coxeter number. $\quad h^{v}=4$.
The Lie algebra $L^{0} . \quad L^{0}=G_{2}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+\mathbb{Z} h_{2} .
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+3 \mathbb{Z} \alpha_{2}
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1,2\right. \\
& \left.2 \alpha_{1}(h)+3 \alpha_{2}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
A^{*}= & \left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1,2\right. \\
& \left.2 \lambda\left(h_{1}\right)+\lambda\left(h_{2}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
& \Phi_{\mathrm{Re}}=\left\{\alpha+r \delta ; \alpha \in \Phi^{0}, r \in \mathbb{Z}\right\} \\
& \Phi_{\mathrm{Im}}=\{k \delta ; k \in \mathbb{Z}, k \neq 0\} \quad \text { Multiplicity } 2 .
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i}, \quad i=1, \ldots, l$ of $L^{0}$.

$$
\omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+2 \gamma, \quad \omega_{2}=\bar{\omega}_{2}+\gamma
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0, \ldots, l \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1, \ldots, l \\
& \langle d, d\rangle=0 .
\end{aligned}
$$

The standard invariant form on $H^{*}$

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0, \ldots, l \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1, \ldots, l \\
& \langle\gamma, \gamma\rangle=0 .
\end{aligned}
$$

## DYNKIN NAME $\tilde{G}_{2}^{\mathrm{t}} \quad$ KAC NAME ${ }^{3} \tilde{D}_{4}$

Dynkin diagram with labelling.


Generalised Cartan matrix.

$$
\begin{aligned}
& 0 \\
& 1 \\
& 2
\end{aligned}\left(\begin{array}{rrr}
0 & 1 & 2 \\
2 & -1 & 0 \\
-1 & 2 & -3 \\
0 & -1 & 2
\end{array}\right)
$$

The integers $a_{0}, a_{1}, \ldots, a_{l}$.


The integers $c_{0}, c_{1}, \ldots, c_{l}$.


The central element $c$.

$$
c=h_{0}+2 h_{1}+3 h_{2}
$$

The basic imaginary root $\delta$.

$$
\delta=\alpha_{0}+2 \alpha_{1}+\alpha_{2}
$$

The element $h_{\theta} \in H_{\mathbb{R}}^{0}$.

$$
h_{\theta}=2 h_{1}+3 h_{2} .
$$

The element $\theta \in\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\theta=2 \alpha_{1}+\alpha_{2}
$$

The Coxeter number. $\quad h=4$.
The dual Coxeter number. $\quad h^{\mathrm{v}}=6$.
The Lie algebra $L^{0} . \quad L^{0}=G_{2}$.
The lattice $M \subset H_{\mathbb{R}}^{0}$.

$$
M=\mathbb{Z} h_{1}+3 \mathbb{Z} h_{2} .
$$

The lattice $M^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
M^{*}=\mathbb{Z} \alpha_{1}+\mathbb{Z} \alpha_{2} .
$$

The fundamental alcove $A \subset H_{\mathbb{R}}^{0}$.

$$
\begin{aligned}
A= & \left\{h \in H_{\mathbb{R}}^{0} ; \alpha_{i}(h)>0 \text { for } i=1,2\right. \\
& \left.2 \alpha_{1}(h)+\alpha_{2}(h)<1\right\} .
\end{aligned}
$$

The fundamental alcove $A^{*} \subset\left(H_{\mathbb{R}}^{0}\right)^{*}$.

$$
\begin{aligned}
& A^{*}=\left\{\lambda \in\left(H_{\mathbb{R}}^{0}\right)^{*} ; \lambda\left(h_{i}\right)>0 \text { for } i=1,2\right. \\
&\left.2 \lambda\left(h_{1}\right)+3 \lambda\left(h_{2}\right)<1\right\} .
\end{aligned}
$$

The root system $\Phi_{0}$ in terms of the root system $\Phi^{0}$ of $L^{0}$.

$$
\begin{aligned}
\Phi_{\mathrm{Re}, \mathrm{~s}}= & \left\{\alpha+r \delta ; \alpha \in \Phi_{\mathrm{s}}^{0}, r \in \mathbb{Z}\right\} \\
\Phi_{\mathrm{Re}, \mathrm{l}}= & \left\{\alpha+3 r \delta ; \alpha \in \Phi_{1}^{0}, r \in \mathbb{Z}\right\} \\
\Phi_{\mathrm{Im}}= & \{3 k \delta ; k \in \mathbb{Z}, k \neq 0\} \quad \text { Multiplicity 2 } \\
& \cup\{(3 k+1) \delta ; k \in \mathbb{Z}\} \quad \text { Multiplicity 1 } \\
& \cup\{(3 k+2) \delta ; k \in \mathbb{Z}\} \quad \text { Multiplicity 1. }
\end{aligned}
$$

The fundamental weights $\omega_{i} \in H_{\mathbb{R}}^{*} \quad i=0, \ldots, l$ in terms of the fundamental weights $\bar{\omega}_{i} \quad i=1, \ldots, l$ of $L^{0}$.

$$
\omega_{0}=\gamma, \quad \omega_{1}=\bar{\omega}_{1}+2 \gamma, \quad \omega_{2}=\bar{\omega}_{2}+3 \gamma .
$$

The standard invariant form on $H$.

$$
\begin{aligned}
& \left\langle h_{i}, h_{j}\right\rangle=a_{j} c_{j}^{-1} A_{i j} \quad i, j=0,1,2 \\
& \left\langle h_{0}, d\right\rangle=1, \quad\left\langle h_{i}, d\right\rangle=0 \quad i=1,2 \\
& \langle d, d\rangle=0 .
\end{aligned}
$$

The standard invariant form on $H^{*}$.

$$
\begin{aligned}
& \left\langle\alpha_{i}, \alpha_{j}\right\rangle=a_{i}^{-1} c_{i} A_{i j} \quad i, j=0,1,2 \\
& \left\langle\alpha_{0}, \gamma\right\rangle=1, \quad\left\langle\alpha_{i}, \gamma\right\rangle=0 \quad i=1,2 \\
& \langle\gamma, \gamma\rangle=0 .
\end{aligned}
$$

## Notation

## Symbol

Meaning
Page of definition

| [xy] | Lie product of elements | 1 |
| :---: | :---: | :---: |
| [HK] | Lie product of subspaces | 1 |
| [A] | the Lie algebra of an associative algebra $A$ | 2 |
| $v \otimes v^{\prime}$ | tensor product | 152 |
| $v \wedge v^{\prime}$ | exterior product | 271 |
| $\langle$, | the Killing form on a finite dimensional Lie algebra | 39 |
| $\langle$, | the standard invariant form on a Kac-Moody algebra | 367 |
| $\langle,\rangle_{t}$ | a bilinear form on the loop algebra | 418 |
| $\langle,\rangle_{0}$ | a contravariant form | 520 |
| \{, \} | a symmetric scalar product | 121 |
| $\succ$ | partial order on weights | 185 |
| [ $V: L(\mu)$ ] | the multiplicity of $L(\mu)$ in $V$ | 459 |
| $a_{0}, a_{1}, \ldots, a_{l}$ | vector associated with an affine Cartan matrix | 386 |
| ad $x$ | the adjoint map | 7 |
| $A=\left(A_{i j}\right)$ | a Cartan matrix | 71 |
| $A=\left(A_{i j}\right)$ | a generalised Cartan matrix (GCM) | 319 |
| $A_{J}$ | a principal minor of $A$ | 344 |
| $A^{0}$ | the underlying Cartan matrix of an affine Cartan matrix $A$ | 394 |
| A | the fundamental alcove | 410 |


| Symbol | Meaning | Page of definition |
| :---: | :---: | :---: |
| $\bar{A}$ | the closure of the fundamental |  |
|  | alcove | 413 |
| $A^{*}$ | the fundamental alcove in the dual |  |
|  | space | 415 |
| $\overline{A^{*}}$ | the closure of the fundamental dual |  |
|  | alcove | 415 |
| $\mathfrak{H}$ | the set of alcoves | 411 |
| $B$ | the subalgebra $H \oplus N$ | 177 |
| $c$ | the Casimir element | 238 |
| $c$ | the canonical central element | 391 |
| $c_{0}, c_{1}, \ldots, c_{l}$ | vector associated with an affine |  |
|  | Cartan matrix | 388 |
| $c(\lambda)$ | scalar action of generalised Casimir |  |
|  | operator | 487 |
| ch $V$ | the character of an $L$-module $V$ | 241 |
| ch $V$ | the character of a module in category $\mathcal{O}$ | 459 |
| C | the fundamental chamber | 112, 247, 378 |
| $\bar{C}$ | closure of the fundamental chamber | 247, 378 |
| $C(V)$ | the Clifford algebra | 282 |
| $C(V)^{+}$ | positive part of the Clifford algebra | 283 |
| $C(V)^{-}$ | negative part of the Clifford algebra | 283 |
| $\mathbb{C}\left[t, t^{-1}\right]$ | the algebra of Laurent polynomials | 417 |
| $d$ | the scaling element | 388 |
| $d_{1}, \ldots, d_{l}$ | degrees of the basic polynomial invariants | 222 |
| $d^{0}(\lambda)$ | the Weyl dimension of an |  |
|  | $L^{0}$-module | 491 |
| $D=\left(d_{i}\right)$ | a diagonal matrix | 110, 390 |
| $D_{0}$ | an endomorphism in the basic |  |
|  | representation | 516 |
| $e_{\alpha}$ | a root vector | 88 |
| $e_{i}$ | a fundamental root vector | 96 |
| $e_{i}$ | a generator of $\tilde{L}(A)$ or $L(A)$ | 323, 332 |
| $e_{\lambda}$ | a characteristic function | 242 |
| $e(\lambda)$ | the characteristic function $e_{\lambda}$ | 487 |
| $E_{i}$ | a generator of $L^{0}$ | 421 |
| $f_{i}$ | a root vector for $-\alpha_{i}$ | 96 |


| Symbol | Meaning | Page of definiti |
| :---: | :---: | :---: |
| $f_{i}$ | a generator of $\tilde{L}(A)$ or $L(A)$ | 323, 332 |
| $F_{i}$ | a generator of $L^{0}$ | 421 |
| $F L(X)$ | the free Lie algebra on a set $X$ | 161 |
| $\mathrm{gl}_{n}(k)$ | the general linear Lie algebra of degree $n$ over $k$ | 5 |
| $G$ | the adjoint group | 207 |
| $h$ | the Coxeter number of a simple Lie algebra | 252 |
| $h$ | the Coxeter number of an affine algebra | 485 |
| $h^{\text {v }}$ | the dual Coxeter number of an affine algebra | 485 |
| $h_{i}$ | a fundamental coroot | 88, 320 |
| $h_{\alpha}$ | the coroot of the root $\alpha$ | 89, 397 |
| $h_{\alpha}^{\prime}$ | the element of $H$ corresponding to $\alpha$ in $H^{*}$ | 46 |
| $h_{\theta}$ | the coroot of $\theta$ | 405 |
| $h_{\alpha}(n)$ | an endomorphism of the basic module | 513 |
| $h(0)$ | an endomorphism of the basic module | 515 |
| ht $\alpha$ | the height of a root $\alpha$ | 62 |
| H | a Cartan subalgebra of a Lie algebra | 23 |
| H | a Cartan subalgebra of a Kac-Moody algebra | 334 |
| $H^{*}$ | the dual space of $H$ | 46 |
| $H_{\mathbb{Q}}$ | a rational vector space in $H$ | 56 |
| $H_{\mathbb{R}}$ | a real vector space in $H$ | 56 |
| $H_{\mathbb{R}}^{*}$ | the dual space of $H_{\mathbb{R}}$ | 57 |
| $H_{i}$ | a generator of $L^{0}$ | 421 |
| $H_{i}$ | a hyperplane | 112 |
| $H_{i}^{+}$ | the positive side of hyperplane $H_{i}$ | 112 |
| ${\underset{\sim}{H}}_{\mathrm{H}_{i}^{-}}$ | the negative side of hyperplane $H_{i}$ | 112 |
| $\tilde{H}$ | the diagonal subalgebra of $\tilde{L}$ | 324 |
| I | the kernel of the map from $\tilde{L}(A)$ to $L(A)$ | 105, 331 |
| $I^{+}$ | the positive subspace of $I$ | 105, 479 |
| $I^{-}$ | the negative subspace of $I$ | 105, 479 |


| Symbol | Meaning | Page of definition |
| :---: | :---: | :---: |
| $J$ | an orbit | 166 |
| $J(\lambda)$ | the maximal submodule of $M(\lambda)$ | 185, 455 |
| $j(\tau)$ | the modular $j$-function | 533 |
| K | a set of positive imaginary roots | 380 |
| $K_{\lambda}$ | the kernel of the map from $\mathfrak{l}(L)$ to $M(\lambda)$ | 178 |
| $K_{A}$ | the set of vectors $u$ with $A u \geq 0$ | 339 |
| $\Omega$ | the generalised partition function | 473 |
| $l$ | the rank of $L$ | 59 |
| $l(w)$ | the length of $w$ | 63 |
| $L$ | a Lie algebra | 1 |
| $L^{n}$ | a power of the Lie algebra $L$ | 7 |
| $L^{(n)}$ | a power of the Lie algebra $L$ | 8 |
| $L_{0, x}$ | the null component of $x$ in $L$ | 23 |
| $L_{\alpha}$ | a root space of $L$ | 36, 333 |
| $L(X, R)$ | the Lie algebra with generators $X$ and relations $R$ | 163 |
| $L(A)$ | the simple Lie algebra with Cartan matrix $A$ | 99 |
| $L(A)$ | the Kac-Moody algebra with GCM A | 331 |
| $L(A)^{\prime}$ | the derived subalgebra of the Kac-Moody algebra $L(A)$ | 335 |
| $\tilde{L}(A)$ | a Lie algebra associated with Cartan matrix $A$ | 99 |
| $\tilde{L}(A)$ | a Lie algebra associated with GCM A | 323, |
| $L(A)^{\sigma}$ | the fixed point subalgebra of $\sigma$ on $L(A)$ | 166 |
| $L(\lambda)$ | the irreducible module with highest weight $\lambda$ | 186, 455, 525 |
| $\tilde{L}_{\alpha}$ | a root space of $\tilde{L}(A)$ | 328 |
| $L^{0}$ | the simple Lie algebra with Cartan matrix $A^{0}$ | 416 |
| $L_{\alpha}$ | a reflecting hyperplane | 246 |
| $L_{\alpha, k}$ | an affine hyperplane | 409 |
| $L_{\theta, 1}$ | a wall of the fundamental alcove | 409 |
| $L_{0}, L_{1}$, | the walls of the fundamental alcove | 412 |


| Symbol | Meaning | Page of definitio |
| :---: | :---: | :---: |
| $\mathfrak{Z}$ | a set of affine hyperplanes | 409 |
| $\mathfrak{R}^{*}$ | a set of affine hyperplanes in the dual space | 414 |
| $\underset{\sim}{\mathcal{Z}}\left(L^{0}\right)$ | the loop algebra of $L^{0}$ | 417 |
| $\tilde{\mathfrak{R}}\left(L^{0}\right)$ | a central extension of the loop algebra | 420 |
| $\hat{\mathfrak{R}}\left(L^{0}\right)$ | realisation of an untwisted Kac-Moody algebra | 420 |
| $\hat{\mathfrak{Z}}\left(L^{0}\right)^{\tau}$ | realisation of a twisted Kac-Moody algebra | 432 |
| $m_{\lambda}$ | a highest weight vector in a Verma module | 180, 452 |
| $m_{\alpha}$ | multiplicity of a root $\alpha$ | 454 |
| M | a lattice | 407 |
| M* | a lattice in the dual space | 413 |
| $M(\lambda)$ | a Verma module | 178, 452 |
| $M^{\perp}$ | the orthogonal subspace of a subspace $M$ | 40 |
| $\mathfrak{M}$ | Monster Lie algebra | 535 |
| $n_{i}$ | an automorphism of $L(A)$ | 373 |
| $n(w)$ | the number of positive roots made negative by $w$ | 63 |
| ${\underset{\sim}{N}}^{\sim}{ }_{\sim}$ | structure constant | 89 |
| $\tilde{N}$ | positive subalgebra of $\tilde{L}(A)$ | 103, 324 |
| $\tilde{N}^{-}$ | negative subalgebra of $\tilde{L}(A)$ | 103, 324 |
| $N$ | positive subalgebra of $L(A)$ | 107, 331 |
| $N^{-}$ | negative subalgebra of $L(A)$ | 107, 331 |
| $N(H)$ | normaliser of a subalgebra $H$ | 23 |
| $\mathcal{O}$ | Bernstein-Gelfand-Gelfand category of modules | 452 |
| $p(k)$ | the number of partitions of $k$ | 507 |
| $p_{l}(k)$ | the number of partitions of $k$ into $l$ colours | 508 |
| $P(L)$ | algebra of polynomial functions on $L$ | 208 |
| $P(L)^{G}$ | G-invariant polynomial functions on $L$ | 210 |


| Symbol | Meaning | Page of definiti |
| :---: | :---: | :---: |
| $P(H){ }^{W}$ | W-invariant polynomial functions |  |
|  | on $H$ | 211 |
| $P S L_{2}(\mathbb{Z})$ | the modular group | 531 |
| $\mathfrak{B}(\lambda)$ | the number of partitions of $\lambda$ into |  |
|  | positive roots | 182 |
| $Q$ | the root lattice | 103, 328 |
| $Q^{+}$ | positive part of the root lattice | 103, 328 |
| $Q^{-}$ | negative part of the root lattice | 103, 328 |
| $Q^{0}$ | root lattice of $L^{0}$ | 404 |
| $Q\left(x_{1}, \ldots, x_{l}\right)$ | quadratic form | 73 |
| $\mathfrak{R}$ | a ring of functions on $H^{*}$ | 242 |
| $s_{i}$ | fundamental reflection | 63, 373 |
| $s_{\alpha}$ | reflection | 60 |
| $s_{\theta}$ | reflection corresponding to root $\theta$ | 405 |
| $s_{\alpha, k}$ | affine reflection | 410 |
| $\mathfrak{E l} n_{n}(\mathbb{C})$ | special linear Lie algebra of degree $n$ over $\mathbb{C}$ | 52 |
| $\operatorname{supp} \alpha$ | support of a root $\alpha$ | 378 |
| Supp $f$ | support of a function $f$ | 241 |
| $S(L)$ | symmetric algebra of $L$ | 201 |
| $S(L)^{G}$ | $G$-invariants in the symmetric algebra of $L$ | 223 |
| $S(H){ }^{W}$ | $W$-invariants in the symmetric algebra of $H$ | 223 |
| $t_{x}$ | a linear map on $H$ | 406 |
| $t_{\alpha}$ | a linear map on $H^{*}$ | 413 |
| $t(M)$ | translation subgroup of the affine |  |
|  | Weyl group | 407 |
| $t\left(M^{*}\right)$ | translation subgroup of the affine |  |
|  | Weyl group | 413 |
| $T(L)$ | tensor algebra of $L$ | 152 |
| $T(V)$ | tensor algebra of $V$ | 324 |
| $T$ | a subalgebra of $L(A)$ | 500 |
| $T^{-}$ | the negative part of $T$ | 500 |
| $\mathfrak{H}(L)$ | universal enveloping algebra of $L$ | 153 |
| $\mathfrak{U}(L)^{+}$ | the ideal $L \mathfrak{H}(L)$ of $\mathfrak{U}(L)$ | 475 |
| $V^{*}$ | the dual module of $V$ | 306 |
| $w_{0}$ | longest element of the Weyl group | 65 |

Symbol
$\left(w_{0}\right)_{J}$
$w_{i}$
$W$
$W$
$W$
$W^{0}$
$W_{J}$
$W^{\sigma}$
$X$
$X^{+}$
$X^{++}$
$Y_{\alpha}(z)$
$Z(L)$
$\alpha_{1}, \ldots, \alpha_{l}$
$\alpha_{1}, \ldots, \alpha_{n}$
$\alpha_{1}, \ldots, \alpha_{n}$
$\alpha^{v}$$\Phi^{-}$

Meaning
longest element of $W_{J} \quad 170$
the weight of $\alpha_{i}$267the Weyl group of a semisimpleLie algebra60the Weyl group of a Kac-Moody
algebra ..... 373
the Weyl group of $L^{0}$ ..... 394
a Weyl subgroup of $W$ ..... 170
the group of $\sigma$-stable elements of $W$ ..... 169
the weight lattice ..... 190, 466
dominant integral weights ..... 190, 466
strictly dominant integral weights ..... 469
a vertex operator ..... 514
the centre of the enveloping algebra ..... 226
fundamental roots of a semisimple Lie algebra ..... 62
fundamental roots of a Kac-Moody algebra ..... 377
dual root ..... 150
the fundamental weight $\omega_{0}$ of an affine algebra ..... 389
component in a vertex operator ..... 515
the basic imaginary root of an affine algebra ..... 384
the Weyl denominator ..... 245
the Kac denominator ..... 469
the Dynkin diagram of a semisimple Lie algebra ..... 80
the Dynkin diagram of a Kac-Moody algebra ..... 353
Dedekind's delta function ..... 533
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the root system of a finite dimensional Lie algebra ..... 36
the root system of a Kac-Moody algebra ..... 377
set of positive roots ..... 58, 377
set of negative roots ..... 58, 377

| Symbol | Meaning | Page of definiti |
| :---: | :---: | :---: |
| $\Phi^{\text {v }}$ | the dual root system of $\Phi$ | 148 |
| $\Phi_{\mathrm{Re}}$ | the real roots in $\Phi$ | 377 |
| $\Phi_{\text {Im }}$ | the imaginary roots in $\Phi$ | 377 |
| $\Phi_{\mathrm{Re}, \mathrm{s}}$ | the short real roots in $\Phi$ | 395 |
| $\Phi_{\text {Re, } 1}$ | the long real roots in $\Phi$ | 395 |
| $\Phi_{\mathrm{Re}, \mathrm{i}}$ | the intermediate real roots in $\Phi$ | 395 |
| $\Phi^{0}$ | the root system of $L^{0}$ | 394 |
| $\Phi_{\text {s }}^{0}$ | the short roots in $\Phi^{0}$ | 400 |
| $\Phi_{1}^{0}$ | the long roots in $\Phi^{0}$ | 400 |
| $\chi_{\lambda}$ | a central character | 226 |
| $\chi^{0}(\lambda)$ | an irreducible character of $L^{0}$ | 487 |
| $\Lambda(V)$ | the exterior algebra of $V$ | 271 |
| $\Lambda^{i}(V)$ | the $i$ th exterior power of $V$ | 271 |
| $\Pi$ | fundamental roots of a semisimple Lie algebra | 58 |
| $\Pi$ | fundamental roots of a Kac-Moody algebra | 320, 334 |
| $\Pi$ | fundamental roots of a Borcherds algebra | 524 |
| $\Pi^{v}$ | fundamental roots in $\Phi^{v}$ | 148 |
| $\Pi^{v}$ | fundamental coroots of a Kac-Moody algebra | 320, 397 |
| $\Pi^{0}$ | fundamental roots in $\Phi^{0}$ | 394 |
| $\Pi_{\mathrm{re}}$ | real fundamental roots of a Borcherds algebra | 525 |
| $\Pi_{i m}$ | imaginary fundamental roots of a Borcherds algebra | 527 |
| $\rho$ | sum of the fundamental weights | 228 |
| $\rho$ | an element satisfying $\rho\left(h_{i}\right)=1$ | 460 |
| $\rho$ | an element in Borcherds' character formula | 526 |
| $\theta$ | the element $\delta-a_{0} \alpha_{0}$ | 404 |
| $\theta_{i}$ | an automorphism of $L$ | 108 |
| $\theta_{1}$ | the highest root | 251 |
| $\theta_{\text {s }}$ | the highest short root | 251 |
| $\theta_{0}$ | an orbit representative | 433 |
| $\omega$ | an automorphism of $L(A)$ | 333 |
| $\tilde{\omega}$ | an automorphism of $\tilde{L}(A)$ | 323 |

Symbol Meaning Page of definition
$\omega_{1}, \ldots, \omega_{l} \quad$ the fundamental weights of a simple Lie algebra ..... 190
$\omega_{0}, \omega_{1}, \ldots, \omega_{l}$ the fundamental weights of an affine Kac-Moody algebra ..... 494
$\Omega$ generalised Casimir operator ..... 461$\Omega_{0} \quad$ an operator on a $T$-module incategory $\mathcal{O}$500

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