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## 1. Definitions and Basic Principles

An algebra over $\mathbb{C}$ is a vector space $A$ over $\mathbb{C}$ with a bilinear multiplication $A \times A \rightarrow A,(x, y) \mapsto x y$. Bilinearity means that for all $x, y \in A$ and $\lambda \in \mathbb{C}$

$$
\begin{aligned}
& \left(x_{1}+x_{2}\right) y=x_{1} y+x_{2} y \\
& x\left(y_{1}+y_{2}\right)=x y_{1}+x y_{2} \\
& (\lambda x) y=x(\lambda y)=\lambda(x y)
\end{aligned}
$$

All the algebras in this course will be over $\mathbb{C}$. The main advantage of $\mathbb{C}$ is that it is algebraically closed.

An associative algebra is an algebra $A$ such that for all $x, y, z \in A, x(y z)=(x y) z$.
A Lie algebra ${ }^{\dagger}$ is an algebra $L$ with multiplication $L \times L \rightarrow L,(x, y) \mapsto[x y]$ such that

$$
\begin{gathered}
{[x x]=0 \text { for all } x \in L,} \\
{[[x y] z]+[[y z] x]+[[z x] y]=0 \text { for all } x, y, z \in L-\text { the Jacobi identity. }}
\end{gathered}
$$

Unless otherwise specified, $L$ shall be an arbitrary Lie algebra. Where dictated by requirements of clarity, we shall write $[x, y]$ for $[x y]$.

Lemma 1.1. For all $x, y \in L,[x y]=-[y x]$.

## Proof.

$$
0=[x+y, x+y]=[x x]+[x y]+[y x]+[y y]=[x y]+[y x]
$$

We say that Lie multiplication is anticommutative.
Lemma 1.2. Suppose $A$ is an associative algebra. Then $A$ can be made into a Lie algebra by defining $[x y]=x y-y x$.

Proof. $[x x]=0$ is clear.

$$
\begin{gathered}
{[[x y] z]=[x y-y x, z]=x y z-y x z-z x y+z y x} \\
{[[y z] x]=y z x-z y x-x y z+x z y} \\
{[[z x] y]=z x y-x z y-y x z+y x z}
\end{gathered}
$$

[^0]All terms cancel, so we have the Jacobi identity.

The Lie algebra obtained in this way is denoted $[A]$.

We can also multiply subspaces in the following way: let $L$ be a Lie algebra and $H, K$ subspaces of $L$. We define $[H K]$ to be the smallest subspace containing all the Lie products $[h k]$ for $h \in H, k \in K$. So

$$
[H K]=\left\{\left[h_{1} k_{1}\right]+\ldots+\left[h_{r} k_{r}\right] \mid h_{i} \in H, k_{i} \in K\right\} .
$$

Lemma 1.3. If $H, K$ are subspaces of $L$ then $[H K]=[K H]$.

## Proof.

$$
[H K]_{\ni}\left[h_{1} k_{1}\right]+\ldots+\left[h_{r} k_{r}\right]=-\left[k_{1} h_{1}\right]-\ldots-\left[k_{r} h_{r}\right] \in[K H]
$$

Multiplication of subspaces is commutative.
$H \subseteq L$ is called a subalgebra of $L$ if $H$ is a subspace of $L$ and $[H H] \subseteq H$. That is, a subalgebra of $L$ is a subset of $L$ that is itself a Lie algebra under the same operations as $L$.

A subset $I \subseteq L$ is called an ideal of $L$ if $I$ is a subspace of $L$ and $[I L] \subseteq I$. We will write $I \triangleleft L$.

Note. Since $[I L]=[L I],[I L] \subseteq I \Leftrightarrow[L I] \subseteq I$.
Every ideal of $L$ is also a subalgebra of $L$, but the converse is not true.
Example. Consider $M_{2}=\{2 \times 2$ matrices over $\mathbb{C}\} .\left[M_{2}\right]$ is a Lie algebra. The subset $T$ of elements of $\left[M_{2}\right]$ of trace zero form an ideal of $\left[M_{2}\right]$. The subset $U$ of elements of $\left[M_{2}\right]$ with upper-right element zero form a subalgebra of $\left[M_{2}\right]$, but not an ideal.

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

Proof. (i) $H \cap K$ is certainly a subspace.

$$
\begin{aligned}
& {[H \cap K, H \cap K] \subseteq[H H] \subseteq H} \\
& {[H \cap K, H \cap K] \subseteq[K K] \subseteq K}
\end{aligned}
$$

(ii)

$$
\begin{aligned}
& {[H \cap K, L] \subseteq[H L] \subseteq H} \\
& {[H \cap K, L] \subseteq[K L] \subseteq K}
\end{aligned}
$$

(iii)

$$
[H+K, H+K] \subseteq[H H]+[H K]+[K H]+[K K] \subseteq H+H+H+K \subseteq H+K
$$

(iv)

$$
[H+K, L] \subseteq[H L]+[K L] \subseteq H+K
$$

Note. The sum of two subalgebras need not be a subalgebra.
We can form factor algebras: let $I \triangleleft L$. In particular, $I$ is an additive subgroup so we can form the factor group $L / I$; the elements of $L / I$ are the cosets $I+x$ for $x \in L$.

$$
\begin{gathered}
(I+x)+(I+y)=I+(x+y) \\
\lambda(I+x)=I+\lambda x
\end{gathered}
$$

We define $[I+x, I+y]=I+[x y]$. We do need to check that this is well-defined, i.e. that if $I+x=I+x^{\prime}$ and $I+y=I+y^{\prime}$ then $I+[x y]=I+\left[x^{\prime} y^{\prime}\right]$. We can find $i_{1}, i_{2} \in I$ such that $x^{\prime}=i_{1}+x$ and $y^{\prime}=i_{2}+y$. So

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

So the coset containing $[x y]$ is the same as that containing $\left[x^{\prime} y^{\prime}\right]$.
It is easy to verify that $L / I$ is a Lie algebra.
A homomorphism of Lie algebras is a linear map $\theta: L_{1} \rightarrow L_{2}$ such that for all $x, y \in L_{1}$, $\theta([x y])=[\theta(x), \theta(y)]$. If $\theta$ is bijective it is called an isomorphism and we write $L_{1} \cong L_{2}$.

Proposition 1.5. Let $\theta: L_{1} \rightarrow L_{2}$ be a homomorphism with kernel $K$. Then $K \triangleleft L_{1}$, $\operatorname{im}(\theta)$ is a subalgebra of $L_{2}$ and $L_{1} / K \cong \operatorname{im}(\theta)$.

Proof. Let $x, y \in L_{1} \cdot[\theta(x), \theta(y)]=\theta[x y] \in \theta\left(L_{1}\right)$, so im $(\theta)$ is a subalgebra of $L_{2}$.
Now let $x \in K$ and $y \in L_{1}$. Then

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

so $y \in K$. Hence, $K \triangleleft L_{1}$.
Now let $x, y \in L_{1}$.

$$
\theta(x)=\theta(y) \Leftrightarrow \theta(x-y)=0 \Leftrightarrow x-y \in K \Leftrightarrow K+x=K+y
$$

So $\theta(x) \mapsto K+x$ is a bijection between $\operatorname{im}(\theta)$ and $L_{1} / K$. We now check that this bijection is an isomorphism of Lie algebras: let $x, y, z \in L_{1}$.

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

So $\operatorname{im}(\theta) \cong L_{1} / K$.

Proposition 1.6. Let $I \triangleleft L$ and $H$ a subalgebra of $L$. Then $I+H$ and $I \cap H$ are subalgebras (by 1.4) and
(i) $I \triangleleft I+H$,
(ii) $I \cap H \triangleleft H$,
(iii) $(I+H) / I \cong H /(I \cap H)$.

Proof. (i) $[I, I+H] \subseteq[I, L] \subseteq I$.
(ii)

$$
\begin{gathered}
{[I \cap H, H]_{\subseteq}[I H]_{\subseteq} \subseteq I} \\
{[I \cap H, H]_{\subseteq}[H H] \subseteq H}
\end{gathered}
$$

(iii) We can form $(I+H) / I$ and $H /(I \cap H)$. Elements of $(I+H) / I$ have the form $I+i+h=I+h$ for $h \in H$. Define a map $\theta: H \rightarrow(I+H) / I$ by $\theta(h)=I+h$. This map is a homomorphism since $\left[I+h, I+h^{\prime}\right]=I+\left[h h^{\prime}\right]$. It is surjective: $\operatorname{im}(\theta)=(I+H) / I$. Consider $\quad \operatorname{ker}(\theta): \quad h \in \operatorname{ker}(\theta) \Leftrightarrow I+h=I \Leftrightarrow h \in I$. So $\quad \operatorname{ker}(\theta)=H \cap I$. By 1.5 $H / \operatorname{ker}(\theta) \cong \operatorname{im}(\theta)$, so $(I+H) / I \cong H /(I \cap H)$.

Note. In this course we shall consider only finite-dimensional Lie algebras over $\mathbb{C}$. In this case,

$$
\operatorname{dim}(L / I)=\operatorname{dim}(L)-\operatorname{dim}(I)
$$

To prove this, select a basis $e_{1}, \ldots, e_{r}$ of $I$ and extend to a basis $e_{1}, \ldots, e_{n}$ of $L$. Each element of $L$ has the form $\lambda_{1} e_{1}+\ldots+\lambda_{n} e_{n}$; each element of $L / I$ has the form $I+\lambda_{r+1} e_{r+1}+\ldots+\lambda_{n} e_{n}=\left(I+\lambda_{r+1} e_{r+1}\right)+\ldots+\left(I+\lambda_{n} e_{n}\right) . I+e_{r+1}, \ldots, I+e_{n}$ form a basis for $L / I$. So $\operatorname{dim}(L / I)=n-r=\operatorname{dim}(L)-\operatorname{dim}(I)$.

Examples. If $\operatorname{dim}(L)=1, L$ has basis $x \cdot[x x]=0$, so $[L L]=0$.
If $\operatorname{dim}(L)=2$ let $x, y$ be a basis for $L .[x x]=[y y]=0$, but $[x y]=-[y x]=$ ? Possibly $[L L]=0$. If $[L L] \neq 0$ then $\operatorname{dim}([L L])=1$. Let $x^{\prime}$ be a basis for $[L L]$ and $x^{\prime}, y^{\prime}$ a basis for $L$ itself. Then we have $\left[x^{\prime} y^{\prime}\right]=\lambda x^{\prime}$ for some $\lambda \in \mathbb{C} \backslash\{0\}$. Re-choose $y^{\prime \prime}=\lambda^{-1} y^{\prime}$ : we then have that $[x y]=x$. Hence, we have two Lie algebras of dimension 2 .

## 2. Representations and Modules of Lie Algebras

Recall that $M_{n}=\{n \times n$ matrices over $\mathbb{C}\}$ and that $\left[M_{n}\right]$ is the Lie algebra of such matrices with $[A B]=A B-B A$.

A representation of a Lie algebra $L$ is a homomorphism $\rho: L \rightarrow\left[M_{n}\right]$ for some $n \in \mathbb{N}$. I.e.,

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

If $\rho$ is a representation of $L$ then so is $\rho^{\prime}$ given by $\rho^{\prime}(x)=T^{-1} \rho(x) T$ where $T$ is a nonsingular $n \times n$ matrix independent of $x$. We say that two representations are equivalent if there is a non-singular $T$ such that $\rho^{\prime}(x)=T^{-1} \rho(x) T$ holds for all $x \in L$.

An $L$-module is a vector space $V$ over $\mathbb{C}$ with a map $V \times L \rightarrow V$ such that (i) $(v, x) \mapsto v x$ is linear in both $v$ and $x$;
(ii) $v[x y]=(v x) y-(v y) x$ for all $v \in V$ and $x, y \in L$.

We shall only deal with finite-dimensional L-modules in this course.
A submodule $W$ of $V$ is a subspace of $V$ such that $w x \in W$ for all $w \in W, x \in L$ i.e. a subspace closed under the right action of the element of $L$.

Proposition 2.1. Let $V$ be an $L$-module with basis $e_{1}, \ldots, e_{n}$. Let $x \in L$ and let $e_{i} x=\sum_{i=1}^{n} \rho_{i j}(x) e_{j}$. Let $\rho(x)$ be the matrix with ij th entry $\rho_{i j}(x)$. Then $x \mapsto \rho(x)$ is a representation of $L$ and a different choice of basis gives an equivalent representation.

Proof. The linear transformation $v \mapsto v x$ has matrix $\rho(x) ; v \mapsto(v x) y$ has matrix $\rho(x) \rho(y) ; \quad v \mapsto(v y) x \quad$ has matrix $\quad \rho(y) \rho(x) ; \quad v \mapsto(v x) y-(v y) x$ has matrix $\rho(x) \rho(y)-\rho(y) \rho(x)$. That is, the linear transformation $v \mapsto v[x y]$ has matrices $\rho[x y]$ and $\rho(x) \rho(y)-\rho(y) \rho(x)$. So $\rho[x y]=\rho(x) \rho(y)-\rho(y) \rho(x)$, so $\rho$ is a representation of $L$.

Now take a new basis $f_{1}, \ldots, f_{n}$ of $V$. The linear transformation $v \mapsto v x$ is represented by a matrix $T^{-1} \rho(x) T$, where $e_{i}=\sum_{i=1}^{n} T_{i j} f_{j}$. So we get a representation $x \mapsto T^{-1} \rho(x) T$ that is equivalent to $\rho$.

An $L$-module is called irreducible if it has no submodules except itself and 0 ; otherwise it is said to be reducible.

An $L$-module $V$ is called decomposable if there are submodules $V_{1}, V_{2} \neq 0$ of $V$ such that $V=V_{1} \oplus V_{2}$; otherwise it is said to be indecomposable.

Proposition 2.2. $L$ is itself an $L$-module under the map $L \times L \rightarrow L:(x, y) \mapsto[x y]$.
This is the adjoint $L$-module; we define ad $y: L \rightarrow L$ by $(\operatorname{ad} y) x=[x y]$.
Proof. It is sufficient to show that for all $x, y, z \in L,[z[x y]]=[[z x] y]-[[z y] x]$. This follows immediately from the Jacobi identity and the anticommutativity of Lie multiplication.

A derivation of a Lie algebra $L$ is a linear map $D: L \rightarrow L$ such that for all $x, y \in L$,

$$
D[x y]=[D x, y]+[x, D y] .
$$

Proposition 2.3. Let $x \in L$. Then ad $x$ is a derivation of $L$.
Proof. Linearity is clear. We need to check that $(\operatorname{ad} x)[y z]=[(\operatorname{ad} x) y, z]+[y,(\operatorname{ad} x) z]$. This is true if and only if the Jacobi identity is true.

Let $V$ be an $L$-module, $W$ a subspace of $V$ and $H$ a subspace of $L$. We define $W H$ to be the subspace spanned by $w h$ for all $w \in W, h \in H$.

If $W$ is a submodule of $V, V / W$ is itself an $L$-module under the action $(W+v) x=W+v x$ for $v \in V, x \in L$. This action is well-defined because

$$
(W+v) x=\left(W+v^{\prime}\right) x \Rightarrow v-v^{\prime} \in W \Rightarrow\left(v-v^{\prime}\right) x \in W \Rightarrow W+v x=W+v^{\prime} x .
$$

## 3. Abelian, Nilpotent and Soluble Lie Algebras

A Lie algebra $L$ is abelian if $[L L]=0$, i.e. $[x y] \equiv 0$.
Define $L^{1}=L$ and inductively define $L^{n+1}=\left[L^{n} L\right]$.
Proposition 3.1. For each $n \in \mathbb{N}, L^{n} \triangleleft L$.
Proof. It is sufficient to show that if $H, K \triangleleft L$ then $[H K] \triangleleft L$. Let $x \in H, y \in K, z \in L$. Is $[[x y] z] \in[H K]$ ?

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

But $y,[y z] \in K$ and $x,[z x] \in H$ so $[H K] L] \subseteq[H K]$.
Clearly $L^{1}=L \triangleleft L$. If we assume inductively that $L^{n} \triangleleft L$ then the above workings show that $L^{n+1}=\left[L^{n} L\right] \triangleleft L$, and the result follows.

Proposition 3.2. $L=L^{1} \supseteq L^{2} \supseteq L^{3} \supseteq \ldots$
Proof. For each $n, L^{n+1}=\left[L^{n} L\right] \subseteq L^{n}$ since $L^{n} \triangleleft L$.
$L$ is nilpotent if there is an $n \in \mathbb{N}$ such that $L^{n}=0$.

Clearly every abelian Lie algebra is nilpotent as $L^{2}=0$.
Example. Let $L$ be the Lie algebra of upper-triangular $n \times n$ matrices with zeroes on the principal diagonal; $\operatorname{dim}(L)=\frac{1}{2} n(n-1) ; L$ is a subalgebra of $\left[M_{n}\right]$. Define subspaces $H_{i}$ by requiring that elements of $H_{i}$ have zeroes on and below the $(i-1)$ th diagonal above the principal diagonal.

Lie multiplication shows that $\left[H_{i}, L\right] \subseteq H_{i-1}$. We show that $L^{i} \subseteq H_{i}$ by induction on $i$. If $i=1$ then $L=L^{1}=H_{1}$. Assume $L^{i} \subseteq H_{i}$ for $i=r$. Then

$$
\begin{aligned}
L^{r+1} & =\left[L^{r} L\right] \\
& \subseteq\left[H^{r} L\right] \\
& \subseteq H_{r+1}
\end{aligned}
$$

In particular, $L^{n} \subseteq H_{n}=0$, so $L$ is nilpotent.

Proposition 3.3. For all $m, n \geq 1,\left[L^{m} L^{n}\right] \subseteq L^{m+n}$.
Proof. Use induction on $n$. If $n=1$ then $\left[L^{m} L^{1}\right]=\left[L^{m} L\right]=L^{m+1}$. Assume for $n=r$ and consider $n=r+1$ :

$$
\begin{aligned}
{\left[L^{m} L^{++1}\right] } & =\left[L ^ { m } \left[L^{\prime} L \rrbracket\right.\right. \\
& \left.=\llbracket L^{\prime} L L^{m}\right] \\
& \left.\left.\left.\left.\subseteq \llbracket L L^{m}\right\rfloor L^{\prime}\right]+\llbracket L^{m} L^{r}\right] L\right] \text { by the Jacobi identity }
\end{aligned}
$$

So

$$
\begin{aligned}
{\left[L^{m} L^{r+1}\right] } & \left.\subseteq\left[L^{m+1} L^{r}\right]+\left[L^{m} L^{r}\right] L\right] \\
& \subseteq L^{m+1+r}+\left[L^{m+r} L\right] \\
& \subseteq L^{m+r+1}
\end{aligned}
$$

We now inductively define another sequence of subspaces of $L$ :

$$
\begin{gathered}
L^{(0)}=L \\
L^{(i+1)}=\left[L^{(i)} L^{(i)}\right]
\end{gathered}
$$

The $L^{(i)}$ are all ideals of $L$, so $L^{(i+1)} \subseteq L^{(i)}$, so $L=L^{(0)} \supseteq L^{(1)} \supseteq L^{(2)} \supseteq \ldots$.
We say that $L$ is soluble if there is an $n \in \mathbb{N}$ such that $L^{(n)}=0$.
Proposition 3.4. (i) $L^{(n)} \subseteq L^{2^{n}}$.
(ii) Every nilpotent Lie algebra is soluble.

Proof. (i) Induction on $n$ : if $n=0 \quad L^{(0)}=L=L^{1}$. Assume for $n=r$ :

$$
\begin{aligned}
L^{(r+1)} & =\left\lfloor L^{(r)} L^{(r)}\right] \\
& \subseteq\left[L^{2^{r}} L^{2^{r}}\right] \\
& \subseteq L^{2^{r+1}}
\end{aligned}
$$

(ii) Suppose $L$ is nilpotent. Then there is an $n$ such that $L^{2^{n}}=0$. By (i) $L^{(n)}=0$ also, so $L$ is soluble.

Example. Let $L$ be the set of all $n \times n$ matrices with zeroes below the principal diagonal. $[L L] \subseteq L ; L$ is a subalgebra of $\left[M_{n}\right]$. Define subspaces $H_{i}$ by requiring that elements of $H_{i}$ have zeroes on and below the $i$ th diagonal above the principal diagonal.

We have that $\left[H_{i} H_{i}\right] \subseteq H_{i+1}$. We show that $L^{(i)} \subseteq H_{i}$ by induction on $i$. If $i=0$ then $L^{(0)}=L=H_{0}$. Assume for $i=r$. Then $L^{(r+1)}=\left[L^{(r)} L^{(r)}\right] \subseteq\left[H_{r} H_{r}\right] \subseteq H_{r+1}$. In particular, $L^{(n)} \subseteq H_{n}$ so $L$ is soluble.

However, $L$ is not nilpotent. To see this, consider

$$
\begin{gathered}
A=\left(\begin{array}{cccc}
0 & \mu_{1} & & 0 \\
& 0 & \ddots & \\
& & \ddots & \mu_{n-1} \\
0 & & 0
\end{array}\right) \text { and } B=\left(\begin{array}{lll}
\lambda_{1} & & 0 \\
& \ddots & \\
0 & & \lambda_{n}
\end{array}\right) \\
A B=\left(\begin{array}{cccc}
0 & \mu_{1} \lambda_{2} & & 0 \\
& 0 & \ddots & \\
& & \ddots & \mu_{n-1} \lambda_{n} \\
0 & & & 0
\end{array}\right) \\
B A=\left(\begin{array}{llll}
0 & \mu_{1} \lambda_{1} & & 0 \\
& 0 & \ddots & \\
& & \ddots & \mu_{n-1} \lambda_{n-1} \\
0 & & & 0
\end{array}\right) \\
{[A B]=A B-B A=\left(\begin{array}{llll}
0 & \mu_{1}\left(\lambda_{2}-\lambda_{1}\right) & & 0 \\
& & 0 & \ddots
\end{array}\right.} \\
0
\end{gathered}
$$

By choosing the $\lambda_{i}$ all unequal and the $\mu_{j}$ suitably we can get any desired matrix

$$
[A B]=\left(\begin{array}{cccc}
0 & v_{1} & & 0 \\
& 0 & \ddots & \\
& & \ddots & v_{n-1} \\
0 & & & 0
\end{array}\right)
$$

Let $K$ be the subspace consisting of all matrices of this form.

$$
K \subseteq[K L] \subseteq[[K L] L] \subseteq \ldots \subseteq[L \ldots L]=L^{i}
$$

So $K \subseteq L^{i}$ for all $i$ and $K \neq 0$, so $L^{i} \neq 0$ for all $i$. Hence, $L$ is not nilpotent.

Proposition 3.5. (i) Every subalgebra of a soluble Lie algebra is soluble.
(ii) Every factor algebra of a soluble Lie algebra is soluble.
(iii) If $I \triangleleft L$ and $I, L / I$ are soluble then $L$ is soluble.

Proof. (i) Let $H$ be a subalgebra of $L . L$ is soluble so $L^{(n)}=0$ for some $n$. We show that $H^{(i)} \subseteq L^{(i)}$ for all $i$. This is true for $i=0$. Assume it is true for $i=r$.

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

So $H^{(n)}=0$ and $H$ is soluble.
(ii) $I \triangleleft L$. We need to show that $(L / I)^{(i)}=\left(I+L^{(i)}\right) / I$. This is true for $i=0$; assume it for $i=r$ :

$$
\begin{aligned}
(L / I)^{(r+1)} & =\left[(L / I)^{(r)},(L / I)^{(r)}\right] \\
& =\left[\frac{I+L^{(r)}}{I}, \frac{I+L^{(r)}}{I}\right] \\
& \left.=\frac{\left[I+L^{(r)}, I+L^{(r)}\right]}{I}\right] \\
& =\frac{I+\left[L^{(r)} L^{(r)}\right]}{I} \\
& =\frac{I+L^{(r+1)}}{I}
\end{aligned}
$$

If $L^{(n)}=0$ then $(L / I)^{(n)}=I / I$, the zero subspace of $L / I$.
(iii) Suppose $I$ and $L / I$ are soluble. $L / I$ soluble $\Leftrightarrow(L / I)^{(m)}=I / I$ for some $m \cdot \frac{I+L^{(m)}}{I}=\frac{I}{I}$, so $L^{(m)} \subseteq I . I$ is soluble so $I^{(n)}=0$ for some $n$.

$$
L^{(m+n)}=\left(L^{(m)}\right)^{(n)} \subseteq L^{(n)}=0
$$

So $L^{(m+n)}=0 ; L$ is soluble.

Proposition 3.6. Let $H, K$ be soluble ideals of $L$. Then $H+K$ is a soluble ideal of $L$.
Proof. We know that $H+K \triangleleft L$. By 1.6, $\frac{H+K}{H} \cong \frac{K}{H \cap K}$. $K$ is soluble, so $K /(H \cap K)$ is soluble. Hence, $(H+K) / H$ is soluble. By the previous proposition, since $H$ is soluble, $H+K$ is soluble.

Corollary 3.7. Any Lie algebra $L$ has a unique maximal soluble ideal.

Proof. Since $\operatorname{dim}(L)<\infty \quad L$ certainly has a maximal soluble ideal. Let $H, K$ be two maximal soluble ideals of $L$. Then $H+K$ is a soluble ideal. $H \subseteq H+K$ and $H$ is maximal, so $H=H+K$. Similarly $K=H+K$, so $H=K$.

This maximal soluble ideal of $L$ is called the soluble radical $L$, usually denoted $R$. If $R=0$ we say that $L$ is semisimple.
$L / R$ is semisimple. For if $R^{\prime} / R$ is the soluble radical of $L / R$ then since $R^{\prime} / R$ and $R$ are both soluble, so is $R^{\prime}$. Hence $R^{\prime} \subseteq R$. So $R^{\prime}=R$ and $R^{\prime} / R=R / R$ is the zero subspace of $L / R$.

A Lie algebra $L$ is called simple if it has no ideals other than 0 and $L .^{\dagger}$
If $\operatorname{dim}(L)=1$ then $L$ is certainly simple. There are other simple Lie algebras. If $L$ is abelian and simple then $\operatorname{dim}(L)=1$, since $[L L]=0$. If $L$ is soluble and simple and $\operatorname{dim}(L)=1$ then $L \neq[L L]$, so $[L L]=0$.

If $I \triangleleft L$ then the ideals of $L / I$ have the form $J / I$ for $J \triangleleft L, J \supseteq I$. So $L / I$ is simple if and only if $I$ is maximal.

A composition series of $L$ is a sequence of subalgebras

$$
L=K_{0} \supset K_{1} \supset \ldots \supset K_{r}=0
$$

where $K_{i+1} \triangleleft K_{i}$ is maximal. The factor algebras $K_{i-1} / K_{i}$ are all simple Lie algebras and are known as the composition factors of $L$.

Proposition 3.8. $L$ is soluble if and only if all composition factors in a composition series of $L$ are 1-dimensional.

Proof. Let $L=K_{0} \supset K_{1} \supset \ldots \supset K_{r}=0$ be a composition series. $L$ is soluble, so $K_{i}$ is soluble, so $K_{i} / K_{i+1}$ is soluble. So $\operatorname{dim}\left(K_{i} / K_{i+1}\right)=1$. Conversely, suppose that $\operatorname{dim}\left(K_{i} / K_{i+1}\right)=1$. Then certainly $K_{i} / K_{i+1}$ is soluble (even abelian). $K_{r-1}$ is soluble and $K_{r-2} / K_{r-1}$ is soluble, so $K_{r-2}$ is soluble. $K_{r-3} / K_{r-2}$ is soluble, so $K_{r-3}$ is soluble. Eventually we see that $K_{0}=L$ is soluble.

[^1]
## 4. Representations of Nilpotent Lie Algebras

We shall first discuss representations of abelian Lie algebras.
Proposition 4.1. Let $L$ be abelian. Then every irreducible $L$-module has dimension 1. Every linear map $L \rightarrow \mathbb{C}$ is a 1-dimensional representation of $L$.

Proof. A 1-dimensional representation of $L$ is by definition a linear map $\lambda: L \rightarrow \mathbb{C}$ such that $\lambda([x y])=\lambda(x) \lambda(y)-\lambda(y) \lambda(x)$. But the RHS of this equation is zero; since $L$ is abelian the LHS is always zero, too. So every linear map $L \rightarrow \mathbb{C}$ is a representation of $L$.

Let $V$ be an irreducible $L$-module and let $x \in L$; consider the linear map $V \rightarrow V$, $v \mapsto v x$. Let $w$ be an eigenvector of this map; i.e. $w \neq 0$ and $w x=\lambda x$ for some $\lambda \in \mathbb{C}$, where $\lambda$ is the eigenvalue. Let $W=\{v \in V \mid v x=\lambda v\}$, the eigenspace. $W$ is a subspace of $V$. Since $w \neq 0, W \neq 0$. We shall show that $W$ is a submodule of $V$.

Let $v \in W, y \in L$.

$$
(v y) x=(v x) y+\underbrace{v[y x]}_{\substack{=0 \\ \because L \text { abelian }}}=(v x) y=(\lambda v) y=\lambda(v y)
$$

So $v y \in W$, which shows that $W$ is a submodule of $V$. But $V$ is irreducible so $V=W$. So $v x=\lambda v$ for all $v \in V$. Hence each $x \in L$ acts on $V$ by scalar multiplication. So every subspace of $V$ is a submodule. Hence $\operatorname{dim}(V)=1$.

We now recall some linear algebra.
Let $A \in M_{n}$. Then the characteristic polynomial of $A$ is $\chi(t)=\operatorname{det}\left(t I_{n}-A\right)$. For nonsingular $T \in M_{n}, A$ and $T^{-1} A T$ have the same characteristic polynomial:

$$
\begin{aligned}
\operatorname{det}\left(t I_{n}-T^{-1} A T\right) & =\operatorname{det}\left(T^{-1}\left(t I_{n}-A\right) T\right) \\
& =\operatorname{det}\left(T^{-1}\right) \operatorname{det}\left(t I_{n}-A\right) \operatorname{det}(T) \\
& =\operatorname{det}\left(t I_{n}-A\right)
\end{aligned}
$$

If $V$ is an $n$-dimensional vector space over $\mathbb{C}$ and $\theta: V \rightarrow V$ is a linear map we define the characteristic polynomial of $\theta$ to be the characteristic polynomial of any matrix representing $\theta$.
$\chi(t) \in \mathbb{C}[t]$ factorizes into linear factors:

$$
\chi(t)=\left(t-\lambda_{1}\right)^{m_{1}} \ldots\left(t-\lambda_{r}\right)^{m_{r}} \text { with } \lambda_{i} \neq \lambda_{j} \text { for } i \neq j
$$

The $\lambda_{i}$ are the eigenvalues, each with multiplicity $m_{i}$.
Question 1: Is there a decomposition of $V$ into a direct sum of subspaces, one for each $\lambda_{i}$ ?

Answer 1: Yes.
There is an eigenvector $v_{i} \in V$ with eigenvalue $\lambda_{i}$, i.e. $\theta\left(v_{i}\right)=\lambda_{i} v_{i}$. The eigenspace for $\theta$ with respect to the eigenvalue $\lambda_{i}$ is

$$
\begin{aligned}
\mathrm{ES}\left(\theta, \lambda_{i}\right) & =\left\{v \in V \mid \theta(v)=\lambda_{i} v\right\} \\
& =\left\{v \in V \mid\left(\theta-\lambda_{i} I\right) v=0\right\}
\end{aligned}
$$

Question 2: $\operatorname{Is} \operatorname{dim}\left(\operatorname{ES}\left(\theta, \lambda_{i}\right)\right)=m_{i}$ ?
Question 3: Is $V$ the direct sum of the eigenspaces of the $\lambda_{i}$ ?
Example. Let $\operatorname{dim}(V)=2$. Let $\left\{e_{1}, e_{2}\right\}$ be a basis for $V$ and take $\theta$ such that $\theta: e_{1} \mapsto e_{2} \mapsto 0$. The matrix of $\theta$ is

$$
\left.\chi(t)=\left\lvert\, \begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right.\right), \left.\begin{array}{cc}
t & -1 \\
0 & t
\end{array} \right\rvert\,=t^{2}
$$

$\theta$ has eigenvalues 0,0 . The eigenspace of $\theta$ with eigenvalue 0 is $\mathbb{C} e_{2}$, so

$$
\operatorname{dim}(\operatorname{ES}(\theta, 0))=1 \neq 2=\text { multiplicity of } 0
$$

So
Answer 2: No.
Answer 3: No.
The generalized eigenspace of $\theta$ with respect to the eigenvalue $\lambda_{i}$ is

$$
\begin{aligned}
\operatorname{GES}\left(\theta, \lambda_{i}\right) & =\left\{v \in V \mid v \text { is annihilated by some power of }\left(\theta-\lambda_{i} I\right)\right\} \\
& =\left\{v \in V \mid \exists N \in \mathbb{N} \text { s.t. }\left(\theta-\lambda_{i} I\right)^{N} v=0\right\}
\end{aligned}
$$

So, in the above example, $\operatorname{GES}(\theta, 0)=V$.
Proposition 4.2. (The Decomposition Theorem) Let $V$ be a vector space of dimension $n$ over $\mathbb{C}$ and let $\theta: V \rightarrow V$ be a linear map with characteristic polynomial

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

with $\lambda_{i}$ distinct and $\sum_{i=1}^{r} m_{i}=n$. Let $V_{i}=\operatorname{GES}\left(\theta, \lambda_{i}\right)$. Then
(i) $V=V_{1} \oplus \ldots \oplus V_{r}$;
(ii) $\operatorname{dim}\left(V_{i}\right)=m_{i}$;
(iii) $\theta\left(V_{i}\right) \subseteq V_{i}$;
(iv) The characteristic polynomial of $\left.\theta\right|_{V_{i}}$ is $\left(t-\lambda_{i}\right)^{m_{i}}$;
(v) $V_{i}=\left\{v \in V \mid\left(\theta-\lambda_{i} I\right)^{m_{i}} v=0\right\}$.

Proof. The proof (omitted) uses the Cayley-Hamilton Theorem, i.e. that $\chi(\theta): V \rightarrow V$ satisfies $\chi(\theta)=0$.

Theorem 4.3. Let $L$ be a nilpotent Lie algebra and $V$ an $L$-module. Let $y \in L$ and $\rho(y): V \rightarrow V: v \mapsto v y$. Then the generalized eigenspaces $V_{i}$ of $V$ with respect to $\rho(y)$ are all submodules of $V$.

Note. This does not hold for arbitrary Lie algebras: we need the nilpotency condition.
Recall. Leibnitz's formula for differentiation:

$$
D^{n}(f g)=\sum_{i=0}^{n}\binom{n}{i}\left(D^{n-i} f\right)\left(D^{i} g\right)
$$

where $D$ denotes the action of differentiation (once).
We first prove
Proposition 4.4. Let $L$ be a Lie algebra and $V$ an $L$-module. Let $v \in V, x, y \in L$ and $\alpha, \beta \in \mathbb{C}$. Then

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

Note. If $\alpha=\beta=0$ then

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

Proof. We use induction on $n$.
If $n=1$ then LHS $=(v x) y-(\alpha+\beta) v x$ and

$$
\begin{aligned}
\mathrm{RHS} & =v(\rho(y)-\alpha I)+v(x(\operatorname{ad} y-\beta I)) \\
& =(v y) x-a v x+v[x y]-\beta v x
\end{aligned}
$$

By the module axioms, LHS $=$ RHS .
Now assume the result for $n=r$.

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

where $x_{i}=x(\operatorname{ad} y-\beta I)^{i}$. Now

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

So

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

which is

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

We shall call the formula of Proposition 4.4 the Leibnitz formula for Lie algebras.
Proof. (of 4.3.) Consider the map $\rho(y): V \rightarrow V$. Let $V_{i}$ be the generalized eigenspace of this map with eigenvalue $\lambda_{i}$. Let $v \in V_{i}$ and $x \in L$. To show that $v x \in V_{i}$ we require that $(v x)\left(\rho(y)-\lambda_{i} I\right)^{N}=0$ for suitably large $N$. Apply Leibnitz with $\alpha=\lambda_{i}, \beta=0$ :

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

$v \in V_{i}$ so $v\left(\rho(y)-\lambda_{i} I\right)^{N-i}=0$ if $N-i$ is sufficiently large. Since $L$ is nilpotent $x(\operatorname{ad} y)^{i}=0$ if $i$ is suitably large. Thus, if $N$ is suitably large, $(v x)\left(\rho(y)-\lambda_{i} I\right)^{N}=0$, and so $v x \in V_{i}$. Thus, each generalized eigenspace is a submodule of $V$.

Corollary 4.5. If $L$ is a nilpotent Lie algebra and $V$ is an indecomposable $L$-module then for all $y \in L$ the linear map $v \mapsto v y$ has only one eigenvalue.

Proof. We know that $V=V_{1} \oplus \ldots \oplus V_{r}$ for generalized eigenspaces $V_{i}$ of $\rho(y)$. These are all submodules. Since $V$ is indecomposable, $r=1$.

Proposition 4.6. Let $L$ be a nilpotent Lie algebra and $V$ an indecomposable $L$-module. Let $y \in L$ have a single eigenvalue $\lambda(y)$ on $V$. Then the map $y \mapsto \lambda(y)$ is a 1dimensional representation of $L$.

Proof. Let $\operatorname{dim}(V)=n$. It is clear that $y \mapsto \lambda(y)$ is linear. We must also show that

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

The RHS is clearly zero, so we need to show that the LHS is zero as well. Consider the trace function:

$$
\begin{aligned}
\operatorname{tr}(A)= & \sum_{i} \alpha_{i i} \\
= & \sum_{\text {eigenvalues of } A} \\
= & -\left(\operatorname{coefficient~of~} t^{n-1} \text { in } \chi(t)\right) \\
& \quad \operatorname{tr}(A B)=\operatorname{tr}(B A)
\end{aligned}
$$

Consider $\rho[x y]: V \rightarrow V$; this has only one eigenvalue, $\lambda[x y]$.

$$
\begin{aligned}
& \operatorname{tr}(\rho[x y])=n \lambda[x y] \\
\operatorname{tr}(\rho[x y])= & \operatorname{tr}(\rho(x) \rho(y)-\rho(y) \rho(x)) \\
= & \operatorname{tr}(\rho(x) \rho(y))-\operatorname{tr}(\rho(y) \rho(x)) \\
= & 0
\end{aligned}
$$

So $n \lambda[x y]=0$ and $\lambda[x y]=0$, as required.

Proposition 4.7. Let $L$ be a nilpotent Lie algebra and $V$ an indecomposable $L$-module. Let $y \in L$ and let $\lambda(y)$ be the unique eigenvalue of $\rho(y)$. Define $\sigma(y): V \rightarrow V$ by

$$
\sigma(y)=\rho(y)-\lambda(y) I .
$$

Then
(i) $\sigma$ is a representation of $L$;
(ii) $\sigma(y)$ is a nilpotent linear map for all $y \in L$.

Proof. (i) We must show that

$$
\begin{aligned}
& \quad \sigma[x y]=\sigma(x) \sigma(y)-\sigma(y) \sigma(x) \\
& \text { RHS }=(\rho(x)-\lambda(x) I)(\rho(y)-\lambda(y) I)-(\rho(y)-\lambda(y) I)(\rho(x)-\lambda(x) I) \\
&=\rho(x) \rho(y)-\rho(y) \rho(x) \\
&=\rho[x y] \\
&=\sigma[x y]+\lambda[x y] I \\
&=\sigma[x y]
\end{aligned}
$$

(ii) $\rho(y)$ has characteristic polynomial $(t-\lambda(y))^{n}$, so $\sigma(y)=\rho(y)-\lambda(y) I$ has characteristic polynomial

$$
\begin{aligned}
\operatorname{det}(t I-\sigma(y)) & =\operatorname{det}((t+\lambda(y)) I-\rho(y)) \\
& =(t+\lambda(y)-\lambda(y))^{n} \\
& =t^{n}
\end{aligned}
$$

So, by the Cayley-Hamilton Theorem, $\sigma(y)$ satisfies $\sigma(y)^{n}=0$.

A representation $\sigma: L \rightarrow\left[M_{n}\right]$ is called a nil representation if each matrix $\sigma(y)$ for $y \in L$ is nilpotent.

Proposition 4.8. Let $L$ be a nilpotent Lie algebra and $\sigma$ a nil representation of $L$. Then $\sigma$ is equivalent to a representation under which each $x \in L$ is represented by a matrix with zeroes on and below the principal diagonal.

Proof. Let $V$ be an $L$-module giving representation $\sigma$. Suppose $V$ is irreducible. $L$ is nilpotent, so $L^{m}=0$ for some $m$. So $V L^{m}=0$. We show that $V L=0$ by descending induction, i.e. that $V L^{i}=0 \Rightarrow V L^{i-1}=0$.

Let $x \in L^{i-1} . \sigma(x)$ is nilpotent so $\sigma(x)^{k}=0$ for some $k$, i.e. $((v x) x) \ldots x=0$ (with $k$ $x$ 's). So there is a $v \in V$ such that $v \neq 0$ and $v x=0$. Let $U$ be the set of all such $v$; we claim $U$ is a submodule of $V$. Let $u \in U, y \in L$.

$$
(u y) x=\underbrace{(u x)}_{=0} y+u[\underbrace{y x]}_{\in L^{i}}=0
$$

So $u y \in U$; hence $U$ is a submodule of $V . U \neq 0$ and $V$ is irreducible, so $U=V$. Hence, $V x=0$ for all $x \in L^{i-1}$, i.e. $V L^{i-1}=0 . V L^{m}=0$ and $V L^{i}=0 \Rightarrow V L^{i-1}=0$, so $V L=0$. But in this situation every subspace of $V$ is a submodule. Since $V$ is irreducible we have that $\operatorname{dim}(V)=1$. So $x \mapsto(0) \in\left[M_{1}\right]$.

If the module $V$ is not irreducible then

$$
V=V_{0} \supset V_{1} \supset \ldots \supset V_{n}=0
$$

where each $V_{i+1}$ is a maximal proper submodule of $V_{i}$. $V$ gives a nil representation, so $V_{i}$ gives a nil representation; $V_{i} / V_{i+1}$ gives a nil representation of $L$. But $V_{i} / V_{i+1}$ is irreducible, so $\operatorname{dim}\left(V_{i} / V_{i+1}\right)=1$ and $\left(V_{i} / V_{i+1}\right) L=0$, i.e. $V_{i} L \subset V_{i+1}$.

Choose a basis $e_{1}, \ldots, e_{n}$ of $V$ adapted to the chain of subspaces, i.e.

$$
\begin{aligned}
& V_{0} \text { has basis } e_{1}, \ldots, e_{n}, \\
& V_{1} \text { has basis } e_{2}, \ldots, e_{n}
\end{aligned}
$$

$V_{i} x \subset V_{i+1}$ so the matrix representing $x$ with respect to this basis has the required upper triangular form.

Corollary 4.9. Let $L$ be a nilpotent Lie algebra and $V$ an indecomposable $L$-module. Then we can choose a basis for $V$ such that the matrix representation of $x$ has the form

$$
\left(\begin{array}{ccc}
\lambda(x) & & * \\
& \ddots & \\
0 & & \lambda(x)
\end{array}\right)
$$

i.e. zeroes below the principal diagonal, and all elements on the principal diagonal equal.

Proof. Follows from 4.5, 4.6, 4.7 and 4.8.

Corollary 4.10. Let $L$ be a nilpotent Lie algebra and $V$ an irreducible $L$-module. Then $\operatorname{dim}(V)=1$.

We now consider arbitrary $L$-modules.
Let $L$ be a nilpotent Lie algebra and $V$ any $L$-module. A weight of $V$ is a 1dimensional representation $\lambda: L \rightarrow \mathbb{C}$ such that there is a $v \in V \backslash\{0\}$ annihilated by some power of $\rho(x)-\lambda(x) I$ for all $x \in L$, where $\rho(x): V \rightarrow V: v \mapsto v x$.

If $\lambda$ is a weight of $V$ the corresponding weight space $V_{\lambda}$ of $V$ is

$$
V_{\lambda}=\{v \in V \mid v \text { annihilated by some power of } \rho(x)-\lambda(x) I \forall x \in L\}
$$

$V_{\lambda}$ is a subspace of $V$.

Theorem 4.11. (The Weight Space Decomposition Theorem) Let $L$ be a nilpotent Lie algebra and $V$ an $L$-module. Then
(i) $V$ has only finitely many weights;
(ii) $V$ is the direct sum of its weight spaces;
(iii) each weight space is a submodule of $V$;
(iv) a basis can be chosen for each $\lambda$-weight space $V_{\lambda}$ such that the matrix representation on $V_{\lambda}$ has the form

$$
x \mapsto\left(\begin{array}{ccc}
\lambda(x) & & * \\
& \ddots & \\
0 & & \lambda(x)
\end{array}\right)
$$

Proof. $V$ may be expressed as a direct sum of indecomposable submodules. Each indecomposable submodule determines a weight $\lambda$ by 4.6. Let $W_{\lambda}$ be the direct sum of all indecomposable components with weight $\lambda$. Then $V=\oplus_{\lambda} W_{\lambda}$. We need to show that $V_{\lambda}=W_{\lambda}$.

Certainly $W_{\lambda} \subseteq V_{\lambda}$. Take $v \in V_{\lambda}$; since $V=\oplus_{\mu} W_{\mu}$ we can write $v=\sum_{\mu} v_{\mu}$ for $v_{\mu} \in W_{\mu}$. For some $N, \quad v(\rho(x)-\lambda(x) I)^{N}=0$. So $\quad \sum_{\mu} v_{\mu}(\rho(x)-\lambda(x) I)^{N}=0 \quad$ and each $v_{\mu}(\rho(x)-\lambda(x) I)^{N} \in W_{\mu}$. So each $v_{\mu}(\rho(x)-\lambda(x) I)^{N}=0$. Suppose $\lambda \neq \mu$, so there is an $x \in L$ such that $\lambda(x) \neq \mu(x)$. By 4.9, $\rho(x)$ is represented on $W_{\mu}$ by a matrix of the form

$$
\left(\begin{array}{ccc}
\mu(x) & & * \\
& \ddots & \\
0 & & \mu(x)
\end{array}\right)
$$

so $\rho(x)-\lambda(x) I$ is represented on $W_{\mu}$ by

$$
\left(\begin{array}{ccc}
\mu(x)-\lambda(x) & & * \\
& \ddots & \\
0 & & \mu(x)-\lambda(x)
\end{array}\right)
$$

Choose $x \in L$ such that $\lambda(x) \neq \mu(x)$. Then the matrix of $\rho(x)-\lambda(x) I$ is non-singular on $W_{\mu}$. So the matrix on $(\rho(x)-\lambda(x) I)^{N}$ on $W_{\mu}$ is non-singular. So

$$
v_{\mu}(\rho(x)-\lambda(x) I)^{N}=0 \Rightarrow v_{\mu}=0 .
$$

So $v_{\mu}=0$ for all $\mu \neq \lambda$. So $v=\sum_{\mu} v_{\mu}=v_{\lambda}$. Hence $v \in W_{\lambda}$, so $V_{\lambda} \subseteq W_{\lambda}$, so $V_{\lambda}=W_{\lambda}$.
(i) $V_{\lambda} \neq 0 \Rightarrow W_{\lambda} \neq 0$ so $\lambda$ is one of the finite number of weights in our decomposition of $V$. So there are only finitely many weights.
(ii) $V=\oplus_{\lambda} W_{\lambda}$ and $V_{\lambda}=W_{\lambda}$ so $V=\oplus_{\lambda} V_{\lambda}$.
(iii) $V_{\lambda}=W_{\lambda}$ is a submodule of $V$.
(iv) Follows from 4.9.

The decomposition $V=\oplus_{\lambda} V_{\lambda}$ is called the weight space decomposition of $V$.

## 5. CARTAN SUBALGEBRAS

Proposition 5.1. Let $L$ be a nilpotent Lie algebra. Then the adjoint representation of $L$ is a nil representation.

Proof. The adjoint representation comes from the $L$-module $L$ itself, ad $x: y \mapsto[y x]$. If $y \in L$ then we have $[y x] \in L^{2},[[y x] x] \in L^{3}$ and so on. But $L$ is nilpotent, so $L^{m}=0$ for some $m \in \mathbb{N}$. I.e., $[[y x] x \ldots x]\left(m-1 x\right.$ 's) is zero. So $(\operatorname{ad} x)^{m-1}=0 ; \operatorname{ad} x$ is nil.

The converse is also true.
Theorem 5.2. (Engel's Theorem) If $L$ is a Lie algebra for which the adjoint representation is a nil representation then $L$ is nilpotent.

Proof. Suppose not. Choose a maximal nilpotent subalgebra $N$ of $L .[L N] \subseteq L$, so we can regard $L$ as an $N$-module. $[N N] \subseteq N$, so $N$ is an $N$-submodule of $L$. Let $M$ be an $N$-submodule of $L$ containing $N$ such that $M / N$ is an irreducible $N$-module. Since $N$ is nilpotent, $\operatorname{dim}(M / N)=1$ by 4.10. So $\operatorname{dim}(M)=\operatorname{dim}(N)+1$.
$L$ gives a nil representation of $L$, and so $L$ gives a nil representation of $N$. So $M$ gives a nil representation of $N$. So $M / N$ gives a nil representation of $N, n \mapsto(\alpha)$, nil if and only if $\alpha=0$. So $(M / N) x \subseteq N / N$ for all $x \in N$. So $[M N] \subseteq N$. Since $\operatorname{dim}(M)=\operatorname{dim}(N)+1, M=N+\mathbb{C} m$. Hence,

$$
\begin{aligned}
{[M M] } & =[N+\mathbb{C} m, N+\mathbb{C} m] \\
& \subseteq[N N]+[N \mathbb{C} m] \\
& \subseteq N+N \\
& =N \\
& \subseteq M
\end{aligned}
$$

Thus, $M$ is a subalgebra of $L$. Since $[N M] \subseteq N$, we also have that $N \triangleleft M$.
We know $M^{2} \subseteq N$. We shall show that for each $i>0$ there exists an integer $n_{i}$ such that $M^{n_{i}} \subseteq N^{i}$. We use induction on $i$. If $i=1$ take $n_{i}=2$, since $M^{2} \subseteq N$. Assume the statement is true for $i=r$. Then we have an $n_{r}$ such that $M^{n_{r}} \subseteq N^{r}$.

$$
\begin{aligned}
M^{n_{r}+1} & =\left[M^{n_{r}} M\right] \\
& =\left[M^{n_{r}}, N+\mathbb{C} m\right] \\
& \subseteq\left[M^{n_{r}} N\right]+M^{n_{r}} \text { ad } m \\
& \subseteq\left[N^{r} N\right]+M^{n_{r}} \text { ad } m \\
& \subseteq N^{r+1}+M^{n_{r}} \text { ad } m
\end{aligned}
$$

We now show by induction on $j$ that $M^{n_{r}+j} \subseteq N^{r+1} M^{n_{r}}(\operatorname{ad} m)^{j}$. This is true for $j=1$; assume it for $j=k$ :

$$
\begin{gathered}
M^{n_{r}+k} \subseteq N^{r+1}+M^{n_{r}}(\operatorname{ad} m)^{k} \\
M^{n_{r}+k+1} \subseteq\left[N^{r+1}, M\right]+\left[M^{n_{r}}(\operatorname{ad} m)^{k}, M\right] \\
N \triangleleft M \Rightarrow N^{r+1} \triangleleft M \Rightarrow\left[N^{r+1}, M\right] \subseteq N^{r+1}
\end{gathered}
$$

So $M^{n_{r}+j} \subseteq N^{r+1} M^{n_{r}}(\operatorname{ad} m)^{j}$.

The adjoint representation of $L$ is nil, so $(\operatorname{ad} m)^{j}=0$ for large $j$. So, for such $j$, $M^{n_{r+1}} \subseteq N^{r+1}$, so there exists an $n_{r}$ such that $M^{n_{r}} \subseteq N^{r}$.
$N$ is nilpotent, so $N^{r}=0$ for some $r$, hence $M^{n_{r}}=0$. So $M$ is nilpotent, which is a contradiction. So $L$ is nilpotent.

We now consider arbitrary Lie algebras. Consider elements $x \in L$ "as far as possible" from 0 , in that ad 0 has all eigenvalues 0 .

We say that $x \in L$ is regular if ad $x: L \rightarrow L$ has as few eigenvalues zero as possible.
Example. Let $L=\left\{A \in\left[M_{2}\right] \mid \operatorname{tr}(A)=0\right\} . \operatorname{dim}(L)=3$. Basis of $L$ :

$$
\begin{gathered}
e=\left(\begin{array}{ll}
0 & 1 \\
0 & 0
\end{array}\right), h=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right), f=\left(\begin{array}{ll}
0 & 0 \\
1 & 0
\end{array}\right) \\
{[h e]=2 e,[h f]=-2 f,[e f]=h}
\end{gathered}
$$

$x \in L$ has the form $a e+b h+c f$ for $a, b, c \in \mathbb{C}$.

$$
\begin{aligned}
& {[e x]=-2 b e+c h} \\
& {[h x]=2 a e-2 c f} \\
& {[f x]=-a h+2 b f}
\end{aligned}
$$

The matrix of ad $x$ with respect to the basis $(e, f, h)$ is

$$
\left(\begin{array}{ccc}
-2 b & c & 0 \\
2 a & 0 & -2 c \\
0 & -a & 2 b
\end{array}\right)
$$

This has characteristic polynomial

$$
\begin{aligned}
\left|\begin{array}{ccc}
t+2 b & -c & 0 \\
-2 a & t & 2 c \\
0 & a & t-2 b
\end{array}\right| & =(t+2 b)\left(t^{2}-2 b t-2 a c\right)+2 a(-c t+2 b c) \\
& =t^{3}-2 b t^{2}-2 a c t+2 b t^{2}-4 b^{2} t-4 a b c-2 a c b+4 a b c \\
& =t^{3}-4 t\left(b^{2}+a c\right)
\end{aligned}
$$

So the multiplicity of zero as an eigenvalue is

$$
1 \text { if } b^{2}+a c \neq 0
$$

$$
3 \text { is } b^{2}+a c=0
$$

So $\left(\begin{array}{cc}a & b \\ c & -a\end{array}\right)$ is regular if and only if $\left|\begin{array}{cc}a & b \\ c & -a\end{array}\right| \neq 0$.
Lemma 5.3. Let $M$ be a subalgebra of $L$. Then the set of all $x \in L$ such that $[M x] \subseteq M$ is a subalgebra $\mathcal{N}(M)$ containing $M$, and $M \triangleleft \mathcal{N}(M)$. Moreover, $\mathcal{N}(M)$ is the largest subalgebra in which $M$ is an ideal.

Proof. Easy - see Exercise Sheet 1.

We call $\mathcal{N}(M)$ the idealizer (or normalizer) of $M$.
Theorem 5.4. Let $x$ be a regular element of $L$. Let $H=\operatorname{GES}(\operatorname{ad} x, 0)$. Then
(i) $H$ is a subalgebra of $L$;
(ii) $H$ is nilpotent;
(iii) $H=\mathcal{N}(H)$.

Proof. (i) Let $y, z \in H$; we need to show that $[y z] \in H$. By Leibnitz,

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

$y \in H$, so $y(\operatorname{ad} x)^{n-i}=0$ if $n-i$ is large. $z \in H$, so $z(\operatorname{ad} x)^{i}=0$ if $i$ is large. Hence, $[y z](\operatorname{ad} x)^{n}=0$ for large $n$. So $[y z] \in H$, and $H$ is a subalgebra.
(iii) We show that $H=\mathcal{N}(H)$. If $z \in \mathcal{N}(H)$ then $[H z] \subseteq H$. Now $x \in H$ since $[x x]=0$. So $[x z],[z x] \in H$. So $[z x]$ is annihilated by some power of ad $x$. So $z \in H$, so $H \supseteq \mathcal{N}(H)$, hence $H=\mathcal{N}(H)$.
(ii) We show that $H$ is nilpotent by Engel's Theorem, i.e. we show that the adjoint representation of $H$ is nil. Let $\operatorname{dim}(L)=n, \operatorname{dim}(H)=l$. Choose a basis $e_{1}, \ldots, e_{l}$ of $H$ and extend to a basis $e_{1}, \ldots, e_{n}$ of $L$. Let $y \in H, y=\lambda_{1} e_{1}+\ldots+\lambda_{l} e_{l}, \lambda_{i} \in \mathbb{C}$. Consider ad $y: L \rightarrow L . H$ is invariant under ad $y$ since $H$ is a subalgebra. Hence we also have $\operatorname{ad} y: H \rightarrow H$ and ad $y: L / H \rightarrow L / H$.

Let $\chi_{L}(t)$ be the characteristic polynomial of ad $y$ on $L$; let $\chi_{H}(t)$ be the characteristic polynomial of ad $y$ on $H$; let $\chi_{L / H}(t)$ be the characteristic polynomial of ad $y$ on $L / H$. We claim that $\chi_{L}(t)=\chi_{H}(t) \chi_{L / H}(t)$. ad $y: L \rightarrow L$ has a matrix of the block form

$$
\begin{aligned}
A & =\left(\begin{array}{cc}
B_{l \times l} & 0_{l \times(n-l)} \\
D_{(n-l) \times l} & C_{(n-l) \times(n-l)}
\end{array}\right) \\
\chi_{L}(t) & =\operatorname{det}\left(t I_{n}-A\right) \\
& =\operatorname{det}\left(\begin{array}{cc}
t I_{l}-B & 0 \\
B & t I_{n-l}-C
\end{array}\right) \\
& =\operatorname{det}\left(t I_{l}-B\right) \operatorname{det}\left(t I_{n-l}-C\right) \\
& =\chi_{H}(t) \chi_{L / H}(t)
\end{aligned}
$$

Given that $y=\lambda_{1} e_{1}+\ldots+\lambda_{l} e_{l}$, how do the coefficients of $\chi_{L}(t), \chi_{H}(t)$ and $\chi_{L / H}(t)$ depend on the $\lambda_{i}$ ? The entries in $A$ are linear functions of the $\lambda_{i}$. The coefficients in $\chi_{L}(t)$ etc. are polynomial functions of the $\lambda_{i}$.

Let $\chi_{L / H}(t)=b_{0}+b_{1} t+b_{2} t^{2}+\ldots$. We claim that $b_{0}$ is not the zero polynomial, for in the special case $y=x \quad b_{0}$ is non-zero. Let $\chi_{H}(t)=t^{m}\left(a_{0}+a_{1} t+a_{2} t^{2}+\ldots\right)$, where $a_{0}$ is not the zero polynomial. We know that $m \leq l$ since $\chi_{H}(t)$ has degree $l$. So $\chi_{L}(t)=t^{m}\left(a_{0} b_{0}+\ldots\right)$ and $a_{0} b_{0}$ is not the zero polynomial. Choose $\lambda_{1}, \ldots, \lambda_{l}$ such that $a_{0} b_{0} \neq 0$. For this $y$ we have that ad $y$ has eigenvalue 0 with multiplicity $m$. So, by the regularity of $x, m \geq l$; hence $m=l$.

So $\chi_{H}(t)=t^{l}\left(a_{0}+\ldots\right)$ has degree $l$, and so is a multiple of $t^{l}$. Hence $\chi_{H}(t)=t^{l}$ since characteristic polynomials are monic. By the Cayley-Hamilton Theorem, ad $y: H \rightarrow H$ satisfies $(\operatorname{ad} y)^{\prime}=0$, so the adjoint representation of $H$ is nil.

Hence, by Engel's Theorem, $H$ is nilpotent.

The generalized eigenspace of ad $x$ with eigenvalue zero where $x \in L$ is regular is called a Cartan subalgebra of $L$.

Any two Cartan subalgebras of $L$ have equal dimension; this is called the rank of $L$.
Any Cartan subalgebra is nilpotent and is its own idealizer.
Example. Let $L=\left\{A \in\left[M_{2}\right] \mid \operatorname{tr}(A)=0\right\}$.

$$
h=\left(\begin{array}{cc}
1 & 0 \\
0 & -1
\end{array}\right)
$$

is a regular element of $L . H$ will be the Cartan subalgebra given by $\operatorname{ad} h . \operatorname{dim}(H)=1$. So

$$
H=\mathbb{C} h=\left\{\left.\left(\begin{array}{cc}
a & 0 \\
0 & -a
\end{array}\right) \right\rvert\, a \in \mathbb{C}\right\}
$$

In general, let $H$ be a Cartan subalgebra of $L$. Then $[L H]_{\subseteq} L$, so we can regard $L$ as an $H$-module and decompose $L$ as $L=\oplus_{\lambda} L_{\lambda}$, where the $L_{\lambda}$ are the weight spaces of $L$ as an $H$-module.

Consider the special case $\lambda=0$, i.e. $0: H \rightarrow \mathbb{C} . L_{0}$ is the 0 -weight space.
Proposition 5.5. $L_{0}=H$. Thus, 0 is a weight of $H$ on $L$.
Proof. By definition,

$$
L_{0}=\left\{y \in L \mid(\operatorname{ad} x)^{k} y=0 \text { for some } k \text { and all } x \in H\right\} .
$$

But $H$ is nilpotent, so $H^{r}=0$, so $[[y x] x \ldots x]$ (with $r-1 x$ 's) is zero. So $H \subseteq L_{0}$.
Now suppose if possible that $H \neq L_{0}$. Then $L_{0} / H$ is an $H$-module. By 4.11 the representation of $H$ on $L_{0} / H$ can be given by matrices

$$
z \mapsto\left(\begin{array}{lll}
0 & & * \\
& \ddots & \\
0 & & 0
\end{array}\right)
$$

So there is a non-zero element of $L_{0} / H$ that is annihilated by all $z \in H$. So there exists a $y \in L_{0} \backslash H$ such that $[y z] \in H$ for all $z \in H$. Hence $y \in \mathcal{N}(H)$. But $H=\mathcal{N}(H)$, so $y \in H$, a contradiction. Hence $H=L_{0}$.

Hence we have the Cartan decomposition of $L$ as

$$
L=H \oplus\left(\oplus_{\lambda \neq 0} L_{\lambda}\right)
$$

The non-zero weights are called roots. Let $\Phi$ be the set of all roots - a finite set. Then

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

By Lie multiplication we know that $[H H]_{\subseteq} H$ and $\left[L_{\alpha} H\right]_{\subseteq} L_{\alpha}$.
Proposition 5.6. Let $\alpha, \beta \in \Phi$. Then
(i) $\left[L_{\alpha} L_{\beta}\right] \subseteq L_{\alpha+\beta}$ if $\alpha+\beta \in \Phi$;
(ii) $\left[L_{\alpha} L_{\beta}\right] \subseteq H=L_{0}$ if $a+\beta=0$;
(iii) $\left[L_{\alpha} L_{\beta}\right]=0$ if $\alpha+\beta \notin \Phi$ and $\alpha+\beta \neq 0$.

Proof. (i) Let $y \in L_{\alpha}, z \in L_{\beta}, x \in H$. By Leibnitz,

$$
[y z](\operatorname{ad} x-(\alpha(x)+\beta(x)) I)^{N}=\sum_{i=0}^{N}\binom{N}{i}\left[y(\operatorname{ad} x-\alpha(x) I)^{N-i}, z(\operatorname{ad} x-\beta(x) I)^{i}\right]
$$

$y \in L_{\alpha}$, so $y$ is annihilated by large powers of $(\operatorname{ad} x-\alpha(x) I)$; similarly $z$ is annihilated by large powers of $(\operatorname{ad} x-\beta(x) I)$. So $[y z](\operatorname{ad} x-(\alpha(x)+\beta(x)) I)^{N}=0$ for large $N$. Hence, $[y z] \in L_{\alpha+\beta}$.
(ii) If $\alpha+\beta=0,[y z] \in L_{0}=H$ by 5.5 , so $\left[L_{\alpha} L_{-\alpha}\right] \subseteq H$.
(iii) If $\alpha+\beta \notin \Phi \cup\{0\}$ then we deduce that $[y z]=0$, otherwise there would be a nonzero element in the $(\alpha+\beta)$-weight space.

Consider $\quad\left[L_{\alpha} L_{-\alpha}\right] \subseteq H \quad$ for $\quad \alpha \in \Phi, \quad$ and $\quad \beta: H \rightarrow \mathbb{C}$. Consider the restriction $\beta:\left[L_{\alpha} L_{-\alpha}\right] \rightarrow \mathbb{C}$.

Proposition 5.7. Let $\alpha \in \Phi$. Consider the subspace $\left[L_{\alpha} L_{-\alpha}\right] \subseteq H$. Let $\beta \in \Phi$. Then $\beta$ restricted to $\left[L_{\alpha} L_{-\alpha}\right]$ is a rational multiple of $\alpha$.

Proof. If $-\alpha \notin \Phi$ then $L_{-\alpha}=0$ and there is nothing to prove, so assume $-\alpha \in \Phi$. Let $\beta \in \Phi$ and consider the functions

$$
\ldots,-2 \alpha+\beta,-\alpha+\beta, \beta, \alpha+\beta, 2 \alpha+\beta, \ldots,
$$

all linear functions on $H$. Since $\Phi$ is finite there exist integers $p, q$ such that

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

are roots but $-(p+1) \alpha+\beta$ and $(q+1) \alpha+\beta$ are not. If $-(p+1) \alpha+\beta=0$ the result is clear; similarly if $(q+1) \alpha+\beta=0$. So we can assume that $-(p+1) \alpha+\beta \neq 0$ and $(q+1) \alpha+\beta \neq 0$.

Let $M=L_{-p \alpha+\beta} \oplus \ldots \oplus L_{q \alpha+\beta} . M$ is a subspace of $L$. Take $y \in L_{\alpha}, z \in L_{-\alpha}$. Then $[y z] \in L_{0}=H$.
$M$ ad $y \subseteq M$ since ad $y$ takes $L_{i \alpha+\beta}$ to $L_{(i+1) \alpha+\beta}$ and $L_{q \alpha+\beta}$ to 0 .
$M$ ad $z \subseteq M$ since ad $y$ takes $L_{i \alpha+\beta}$ to $L_{(i-1) \alpha+\beta}$ and $L_{-p \alpha+\beta}$ to 0 .
Let $x=[y z] \in H ; M$ ad $x \subseteq M$ by the above. We now calculate $\operatorname{tr}_{M}(\operatorname{ad} x)$ in two different ways:

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

$\operatorname{ad} x$ acts on $L_{i \alpha+\beta}$ as

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

by 4.11 . Hence,

$$
\operatorname{tr}_{L_{i \alpha+\beta}}(\operatorname{ad} x)=(i a+\beta)(x) \operatorname{dim}\left(L_{i \alpha+\beta}\right)
$$

So

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

Equating the two traces gives

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

And so

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

So there exists an $r_{\beta, \alpha} \in \mathbb{Q}$ such that $\beta=r_{\beta, \alpha} \alpha$ on $\left[L_{\alpha} L_{-\alpha}\right]$.

## 6. The Killing Form

We define a map $L \times L \rightarrow \mathbb{C}$ by $(x, y) \mapsto \operatorname{tr}(\operatorname{ad} x \operatorname{ad} y)$. Define $\langle x, y\rangle=\operatorname{tr}(\operatorname{ad} x \operatorname{ad} y)$. The map $(x, y) \mapsto\langle x, y\rangle$ is called the Killing form.

Proposition 6.1. (i) $\langle\cdot$,$\rangle is bilinear;$
(ii) $\langle\cdot, \cdot\rangle$ is symmetric;
(iii) $\langle\cdot, \cdot\rangle$ is invariant, i.e. $\langle[x y], z\rangle=\langle x,[y z]\rangle$ for all $x, y, z \in L$.

Proof. (i) Easy.
(ii) Follows from the identity $\operatorname{tr}(A B)=\operatorname{tr}(B A)$.
(iii)

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

The Killing form is called non-degenerate if $\langle x, y\rangle=0 \quad \forall y \in L \Rightarrow x=0$.

The Killing form is identically zero if $\langle x, y\rangle=0 \quad \forall x, y \in L$.

Proposition 6.2. Let $I \triangleleft L$ 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. Choose a basis of $I$ and extend to a basis of $L$. With respect to this basis, since $I$ is an ideal, ad $x$ is represented by a matrix of the block form

$$
\left(\begin{array}{ll}
A_{1} & 0 \\
A_{2} & 0
\end{array}\right)
$$

Similarly, ad $y$ is represented by

$$
\left(\begin{array}{ll}
B_{1} & 0 \\
B_{2} & 0
\end{array}\right)
$$

So ad $x \operatorname{ad} y$ is represented by

$$
\left(\begin{array}{ll}
A_{1} B_{1} & 0 \\
A_{2} B_{1} & 0
\end{array}\right)
$$

Thus, $\operatorname{tr}(\operatorname{ad} x \operatorname{ad} y)=\operatorname{tr}\left(A_{1} B_{1}\right)$. But $A_{1}$ is the matrix of ad $x$ on $I$ and $B_{1}$ is the matrix of ad $y$ on $I$. Thus,

$$
\langle x, y\rangle_{I}=\operatorname{tr}\left(A_{1} B_{1}\right)=\langle x, y\rangle_{L} .
$$

For any subspace $M$ of $L$ define the perpendicular space $M^{\perp}$ by

$$
M^{\perp}=\{x \in L \mid\langle x, y\rangle=0 \forall y \in M\} .
$$

$M^{\perp}$ is also a subspace of $L$.
Lemma 6.3. $I \triangleleft L \Rightarrow I^{\perp} \triangleleft L$.
Proof. Let $x \in I^{\perp}$ and $y \in L$; we show that $[x y] \in I^{\perp}$. Let $z \in I$.

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

$[y z] \in I, x \in I^{\perp}$. Hence $[x y] \in I^{\perp}$.

In particular, $L^{\perp} \triangleleft L$ :

$$
L^{\perp}=\{x \in L \mid\langle x, y\rangle=0 \forall y \in L\}
$$

So $L^{\perp}=0 \operatorname{iff}\langle\cdot, \cdot\rangle$ is non-degenerate; $L^{\perp}=L$ iff $\langle\cdot, \cdot\rangle$ is identically zero.
Proposition 6.4. Let $L$ be a Lie algebra with $L \neq 0, L^{2}=L$. Let $H$ be a Cartan subalgebra of $L$. Then there is an $x \in H$ such that $\langle x, x\rangle \neq 0$.

Proof. Consider the Cartan decomposition of $L$ as an $H$-module:

$$
\begin{gathered}
L=\oplus_{\lambda} L_{\lambda} \\
L^{2}=[L L]=\left[\oplus_{\lambda} L_{\lambda}, \oplus_{\mu} L_{\mu}\right]=\sum_{\lambda, \mu}\left[L_{\lambda} L_{\mu}\right]
\end{gathered}
$$

Now $\left.\mid L_{\lambda} L_{\mu}\right\rfloor \subseteq L_{\lambda+\mu}$ by 5.6 , where $L_{\lambda+\mu}=0$ if $\lambda+\mu$ is not a weight. Now $L_{0}=H$, so, since $L^{2}=L$, we have

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

Now $L$ is not nilpotent since $L^{2}=L$, but $H$ is nilpotent, so $H \neq L$. So there exists a root $\beta \in \Phi . \beta$ is a 1-dimensional representation of $H, \beta \neq 0 . \beta$ vanishes on $[H H]$ since if $x, y \in H$ then $\beta[x y]=\beta(x) \beta(y)-\beta(y) \beta(x)=0$. So $H=[H H]$.

Hence, there exists an $\alpha \in \Phi$ with $\left[L_{\alpha} L_{-a}\right] \neq 0$. In particular, $L_{-\alpha} \neq 0$. Also, $\beta$ does not vanish on $\left[L_{\alpha} L_{-\alpha}\right]$. Choose $x \in\left[L_{\alpha} L_{-\alpha}\right]$ such that $\beta(x) \neq 0$. By definition, $\langle x, x\rangle=\operatorname{tr}(\operatorname{ad} x \operatorname{ad} x) . \operatorname{ad} x$ acts on $L_{\lambda}$ by

$$
\left(\begin{array}{ccc}
\lambda(x) & & 0 \\
& \ddots & \\
0 & & \lambda(x)
\end{array}\right)
$$

by 4.11. So $(\operatorname{ad} x)^{2}$ acts on $L_{\lambda}$ by

$$
\left(\begin{array}{ccc}
\lambda(x)^{2} & & 0 \\
& \ddots & \\
0 & & \lambda(x)^{2}
\end{array}\right)
$$

So $\langle x, x\rangle=\sum_{\lambda} \operatorname{dim}\left(L_{\lambda}\right) \lambda(x)^{2}$. However, by 5.7 there exists an $r_{\lambda, \alpha} \in \mathbb{Q}$ such that $\lambda(x)=r_{\lambda, \alpha} \alpha(x)$ since $x \in\left[L_{\alpha} L_{-\alpha}\right]$. So $\langle x, x\rangle=\alpha(x)^{2} \sum_{\lambda} \operatorname{dim}\left(L_{\lambda}\right) r_{\lambda, \alpha}^{2}$. In particular, $\beta(x)=r_{\beta, \alpha} \alpha(x)$. Now $\beta(x) \neq 0$, so $\alpha(x)=0$ and $r_{\beta, \alpha} \neq 0$. So

$$
\langle x, x\rangle=\underbrace{\alpha(x)^{2}}_{\neq 0} \underbrace{\sum_{\lambda} \operatorname{dim}\left(L_{\lambda}\right) r_{\lambda, \alpha}^{2}}_{>0}
$$

So $\langle x, x\rangle \neq 0$.

Theorem 6.5. If the Killing form of $L$ is identically zero then $L$ is soluble.

Proof. Use induction on $\operatorname{dim}(L)$. If $\operatorname{dim}(L)=1$ then $L$ is certainly soluble. If $\operatorname{dim}(L)>1$, $L^{2} \neq L$ by 6.4, and $L^{2} \triangleleft L$. The Killing form on $L^{2}$ is the restriction of that on $L$, and so is identically zero. $\operatorname{dim}\left(L^{2}\right)<\operatorname{dim}(L)$; by induction $L^{2}$ is soluble. $L / L^{2}$ is abelian, and so soluble. Hence $L$ is soluble.

Note. The converse is not true. Consider a Lie algebra of dimension 2, basis $\{x, y\}$, with $[x y]=x$.

$$
\operatorname{ad} y \mapsto\left(\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right),(\operatorname{ad} y)^{2} \mapsto\left(\begin{array}{ll}
1 & 0 \\
0 & 0
\end{array}\right)
$$

So $\langle y, y\rangle=\operatorname{tr}\left((\operatorname{ad} y)^{2}\right)=1$.

For which Lie algebras is the Killing form non-degenerate?
Let $R$ be the soluble radical of $L$. Then $L$ is semisimple if and only if $R=0$.
Theorem 6.6. The Killing form on $L$ is non-degenerate if and only if $L$ is semisimple.
Proof. Suppose the Killing form on $L$ is degenerate. Then $L^{\perp} \neq 0 . L^{\perp} \triangleleft L$ by 6.3. So the Killing form of $L$ restricted to $L^{\perp}$ is the Killing form of $L^{\perp}$. Hence, the Killing form of $L^{\perp}$ is identically zero, since $x, y \in L^{\perp} \Rightarrow\langle x, y\rangle=0$. Hence, by $6.5, L^{\perp}$ is soluble. $L^{\perp}$ is a non-zero soluble ideal of $L$, and so $L$ is not semisimple.

Conversely, suppose $L$ is not semisimple. Then $R \neq 0$, so

$$
R=R^{(0)} \supseteq R^{(1)} \supseteq \ldots R^{(k-1)} \supseteq R^{(k)}=0
$$

Let $I=R^{(k-1)} ; I \neq 0$ but $I^{2}=0$, and $I \triangleleft L$. Choose a basis of $I$ and extend to a basis of $L$. Let $x \in I$ and $y \in L$. With respect to this basis, ad $x: L \rightarrow L$ has matrix

$$
\left(\begin{array}{ll}
0 & 0 \\
A & 0
\end{array}\right)
$$

ad $y: L \rightarrow L$ has matrix

$$
\left(\begin{array}{cc}
B_{1} & 0 \\
B_{2} & B_{3}
\end{array}\right)
$$

So ad $x$ ad $y: L \rightarrow L$ has matrix

$$
\left(\begin{array}{cc}
0 & 0 \\
A B_{1} & 0
\end{array}\right)
$$

Hence $\langle x, y\rangle=\operatorname{tr}(\operatorname{ad} x \operatorname{ad} y)=0$. Thus $I \subseteq L^{\perp}$. Thus $L^{\perp} \neq 0$, which implies that the Killing form is degenerate.

Suppose $L_{1}, L_{2}$ are Lie algebras. We can define the direct sum $L_{1} \oplus L_{2}$ to be the set of pairs $\left(x_{1}, x_{2}\right) \in L_{1} \times L_{2}$ with $\left[\left(x_{1}, x_{2}\right)\left(y_{1}, y_{2}\right)\right]=\left(\left[x_{1} y_{1}\right],\left[x_{2} y_{2}\right]\right)$. Similarly, we can define $L_{1} \oplus \ldots \oplus L_{k}$.

The 1-dimensional Lie algebra is simple, and is called the trivial simple Lie algebra.
Theorem 6.7. A Lie algebra is semisimple if and only if it is the direct sum of simple nontrivial Lie algebras.

Proof. Let $L$ be a semisimple Lie algebra. If $L$ is simple $L$ is nontrivial and there is nothing to prove. So assume $L$ is not simple. Choose a minimal non-zero ideal $I \triangleleft L$, $I \neq L$. Consider $I^{\perp}: I^{\perp} \triangleleft L$ as well.

$$
x \in I^{\perp} \Leftrightarrow\langle x, y\rangle=0 \text { for all } y \in I
$$

The Killing form is non-degenerate on $L$. This gives $\operatorname{dim}(I)$ linearly independent conditions on $x$, since the form is non-degenerate. So $\operatorname{dim}\left(I^{\perp}\right)=\operatorname{dim}(L)-\operatorname{dim}(I)$ by the Rank-Nullity Formula.

Now consider $I \cap I^{\perp} . I \cap I^{\perp}$ is an ideal, so the Killing form on $I \cap I^{\perp}$ is the restriction of that on $L$. But $x, y \in I \cap I^{\perp} \Rightarrow\langle x, y\rangle=0$. So $I \cap I^{\perp}$ is soluble by 6.5. Since $L$ is semisimple, $I \cap I^{\perp}=0$.

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

So $L=I+I^{\perp}$, and $I \cap I^{\perp}=0$, so $L=I \oplus I^{\perp}$ as a direct sum of subspaces. Let $x \in I$ and $y \in I^{\perp}$. Then $[x y] \in\left[I I^{\perp}\right] \subseteq I \cap I^{\perp}=0$. So $L=I \oplus I^{\perp}$ as a direct sum of Lie algebras:

$$
\left[a+b, a^{\prime}+b^{\prime}\right]=\left[a a^{\prime}\right]+\left[b b^{\prime}\right]
$$

We now show that $I$ is simple. Let $J \triangleleft I$.

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

So $J$ is an ideal of $L$ contained in $I$. But $I$ is minimal, so either $J=0$ or $J=I$. Hence, $I$ is simple.

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

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

Thus, $J$ is a soluble ideal of $L$. Since $L$ is semisimple, $J=0$, and so $I^{\perp}$ is semisimple.
So $L=I \oplus I^{\perp}$, $I$ simple and $I^{\perp}$ semisimple. By induction, $I^{\perp}$ is the direct sum of simple nontrivial Lie algebras. Hence, $L$ has this property.

Conversely, let $L=L_{1} \oplus \ldots \oplus L_{r}$, where each $L_{i}$ is simple and nontrivial. $L_{i}$ is semisimple, and so its Killing form is non-degenerate by 6.6. For each $i, L_{i} \triangleleft L$. If $i \neq j$ and $x_{i} \in L_{i}, \quad x_{j} \in L_{j}$, then $\left\langle x_{i}, x_{j}\right\rangle=0$. For if $y \in L$ then $y \mathrm{ad} x_{i}=\left[y x_{i}\right] \in L_{i}$; $y$ ad $x_{j}=\left[y x_{j}\right] \in L_{j}$. So $y$ ad $x_{i}$ ad $x_{j}=\left[\left[y x_{i}\right] x_{j}\right] \in L_{i} \cap L_{j}=0$. Thus $\operatorname{ad} x_{i}$ ad $x_{j}=0$, so $\operatorname{tr}\left(\operatorname{ad} x_{i} \operatorname{ad} x_{j}\right)=0$, so $\left\langle x_{i}, x_{j}\right\rangle=0$.

We show $L^{\perp}=0$. Let $x \in L^{\perp}, \quad x=x_{1}+\ldots+x_{r}, \quad x_{i} \in L_{i}$. Let $y_{i} \in L_{i}$. Then $\left\langle x, y_{i}\right\rangle=\left\langle x_{i}, y_{i}\right\rangle=0$, since $x \in L^{\perp}$. So $\left\langle x_{i}, y_{i}\right\rangle=0$ for all $y_{i} \in L_{i}$. This implies $x_{i}=0$ since the Killing form on $L_{i}$ is non-degenerate. So $x=0$, and so $L^{\perp}=0$. Hence the Killing form on $L$ is non-degenerate, and so $L$ is semisimple.

## 7. The Lie Algebra $\mathfrak{s l}_{n}(\mathbb{C})$

$$
\mathfrak{s l}_{n}(\mathbb{C})=\left\{A \in\left[M_{n}\right] \mid \operatorname{tr}(A)=0\right\}
$$

$\mathfrak{s l}_{n}(\mathbb{C})$ is an ideal of the Lie algebra $\left[M_{n}\right]=\mathfrak{g l}_{n}(\mathbb{C})$.

Theorem 7.1. $\mathfrak{s l}_{n}(\mathbb{C})$ is a simple Lie algebra.
Proof. Every ideal of $\mathfrak{s l}_{n}(\mathbb{C})$ is an ideal of $\left[M_{n}\right]$. We shall show that if $I \triangleleft\left[M_{n}\right]$ and $I \subseteq \mathfrak{s l}_{n}(\mathbb{C})$ then $I=\mathfrak{s l}_{n}(\mathbb{C})$ or $I=0$.

Let $I$ be a non-zero ideal of $\left[M_{n}\right]$ contained in $\mathfrak{s l}_{n}(\mathbb{C})$. Let $x \in I \backslash\{0\}$. Then

$$
x \in \sum_{p, q} x_{p q} E_{p q}, x_{p q} \in \mathbb{C} \text { not all zero, }
$$

where $E_{p q}$ is an $n \times n$ matrix with 1 in the $(p, q)$ th position and zeroes elsewhere.
Case 1: Suppose $\exists i \neq j$ such that $x_{i j} \neq 0$. Then

$$
\begin{gathered}
{\left[E_{i i} x\right]=\sum_{q} x_{i q} E_{i q}-\sum_{p} x_{p i} E_{p i} \in I} \\
{\left[\left[E_{i i} x\right] E_{j j}\right]=x_{i j} E_{i j}+x_{j i} E_{j i} \in I} \\
{\left[E_{i i}-E_{j j}, x_{i j} E_{i j}-x_{j i} E_{j i}\right]=2 x_{i j} E_{i j}-2 x_{j i} E_{j i} \in I}
\end{gathered}
$$

So $4 x_{i j} E_{i j} \in I, x_{i j} \neq 0$, so $E_{i j} \in I$.
Case 2: Suppose $x_{i j}=0 \forall i \neq j$. Then $x=\sum_{p} x_{p p} E_{p p} ; \sum_{p} x_{p p}=0$, so not all the $x_{p p}$ are equal. Suppose $x_{i i} \neq x_{i j}$.

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

So $E_{i j} \in I$.

So in either case $E_{i j} \in I$ for some $i \neq j$. Let $q \neq i, j$. Then $\left[E_{i j}, E_{i q}\right]=E_{i q} \in I$. So $I$ contains all $E_{i q}$ with $q \neq i$. Let $p \neq i, q$. Then $\left[E_{p i}, E_{i q}\right]=E_{p q} \in I$. So $E_{p q} \in I$ for all $p \neq q$. For $p \neq q,\left[E_{p q}, E_{q p}\right]=E_{p p}-E_{q q} \in I$. So $I=\mathfrak{s l}_{n}$. Hence, $\mathfrak{s l}_{n}(\mathbb{C})$ is simple.

It is easy to see that $\operatorname{dim}\left(\mathfrak{s l}_{n}(\mathbb{C})\right)=n^{2}-1$; assume $n \geq 2$. We now find a regular element of $\mathfrak{s l}_{n}(\mathbb{C})$.

Proposition 7.2. Let

$$
x=\left(\begin{array}{lll}
\lambda_{1} & & 0 \\
& \ddots & \\
0 & & \lambda_{n}
\end{array}\right) \in \mathfrak{s l}_{n}(\mathbb{C})
$$

with $\sum_{i} \lambda_{i}=0$ and $i \neq j \Rightarrow \lambda_{i} \neq \lambda_{j}$. Then $x$ is regular.
Proof. Take a basis of $\mathfrak{s l}_{n}(\mathbb{C})$ :

$$
\left\{E_{i j} \mid i \neq j\right\} \cup\left\{E_{11}-E_{22}, E_{22}-E_{33}, \ldots, E_{n-1, n-1}-E_{n n}\right\}
$$

Consider the matrix of ad $x$ with respect to this basis:

$$
\begin{gathered}
{\left[E_{i j} x\right]=\left(\lambda_{j}-\lambda_{i}\right) E_{i j} \text { for } i \neq j} \\
{\left[E_{i i}-E_{i+1, i+1}, x\right]=0}
\end{gathered}
$$

So the $\left(n^{2}-1\right) \times\left(n^{2}-1\right)$ matrix of ad $x$ with respect to this basis is

$$
\left(\begin{array}{lllll}
\ddots & & & & 0 \\
& \lambda_{j}-\lambda_{i} & & & \\
& & \ddots & & \\
& & & 0 & \\
0 & & & & \ddots
\end{array}\right)
$$

The characteristic polynomial is $t^{n-1} \prod_{i \neq j}\left(t-\lambda_{i}+\lambda_{j}\right)$; the multiplicity of zero as an eigenvalue of ad $x$ is $n-1$.

Now let $y \in \mathfrak{s l}_{n}(\mathbb{C})$. ad $y$ is similar to a matrix in Jordan canonical form, say

$$
\left(\begin{array}{ccc}
J_{m_{1}}\left(\mu_{1}\right) & & 0  \tag{*}\\
& \ddots & \\
0 & & J_{m_{r}}\left(\mu_{r}\right)
\end{array}\right)
$$

where $J_{m}(\mu)$ is the $m \times m$ matrix

$$
\left(\begin{array}{llll}
\mu & 1 & & 0 \\
& \mu & \ddots & \\
& & \ddots & 1 \\
0 & & & \mu
\end{array}\right)
$$

$J_{m}(\mu)=\mu I_{m}+J_{m}$, where $J_{m}=J_{m}(0) . J_{m}^{k}$ has 1's on the $k$ th diagonal above the principal diagonal and zeroes elsewhere. We claim that $J_{m}(\mu)$ commutes with any matrix of the form $\alpha I_{m}+\beta J_{m}+\beta J_{m}^{2}+\ldots$, so $(*)$ commutes with all matrices of the block form

$$
\left(\begin{array}{ccccccccc}
\alpha_{1} & \beta_{1} & & 0 & & & & & \\
& \alpha_{1} & \ddots & & & & & & \\
& & \ddots & \beta_{1} & \cdots & & & 0 & \\
0 & & & \alpha_{1} & & & & & \\
& & \vdots & & & \ddots & & & \vdots \\
& & \\
& & & & & \alpha_{r} & \beta_{r} & & 0 \\
& 0 & & & \ldots & & \alpha_{r} & \ddots & \\
& & & & & & & \ddots & \beta_{r} \\
& & & & & 0 & & & \alpha_{r}
\end{array}\right)
$$

These matrices form a vector space of dimension $m_{1}+m_{2}+\ldots+m_{r}=n$.
Assume $m_{1} \alpha_{1}+\ldots+m_{r} \alpha_{r}=0$, i.e. the matrix is in $\mathfrak{s l}_{n}(\mathbb{C})$. We have a vector space of dimension $n-1$. All these matrices lie in the zero eigenspace of ad $y$, so the multiplicity of zero as an eigenvalue of ad $y$ is at least $n-1$.

Thus, $x$ is regular.

Proposition 7.3. The subalgebra of diagonal matrices in $\mathfrak{s l}_{n}(\mathbb{C})$ is a Cartan subalgebra.
Proof. Let

$$
x=\left(\begin{array}{lll}
\lambda_{1} & & 0 \\
& \ddots & \\
0 & & \lambda_{n}
\end{array}\right)
$$

$\lambda_{i} \neq \lambda_{j}$ for $i \neq j, \sum_{i} \lambda_{i}=0 . x$ is regular. Let

$$
H=\left\{\left(\begin{array}{lll}
* & & 0 \\
& \ddots & \\
0 & & *
\end{array}\right) \in \operatorname{sl}_{n}(\mathbb{C})\right\}
$$

$\operatorname{dim}(H)=n-1$. If $y \in H$ then $[y x]=0$, so $H \subseteq \operatorname{ES}(\operatorname{ad} x, 0) \subseteq \operatorname{GES}(\operatorname{ad} x, 0)$. But

$$
\operatorname{dim}(H)=n-1=\operatorname{dim}(\operatorname{GES}(\operatorname{ad} x, 0)) .
$$

Thus, $H$ is a Cartan subalgebra of $\mathfrak{s l}_{n}(\mathbb{C})$.

Proposition 7.4. Let $L=\mathfrak{s l}_{n}(\mathbb{C})$ and let $H$ be the diagonal subalgebra. Then

$$
L=H \oplus\left(\oplus_{i \neq j} \mathbb{C} E_{i j}\right)
$$

is the Cartan decomposition of $L$ with respect to $H$.
Proof. Clearly $L=H \oplus\left(\oplus_{i \neq j} \mathbb{C} E_{i j}\right)$ as a direct sum of vector spaces. Let

$$
h=\lambda_{1} E_{11}+\ldots+\lambda_{n} E_{n n} \in H .
$$

Then $\left\lfloor E_{i j} h\right\rfloor=\left(\lambda_{j}-\lambda_{i}\right) E_{i j}$. So $\mathbb{C} E_{i j}$ is an $H$-module giving a 1-dimensional representation

$$
\left(\begin{array}{ccc}
\lambda_{1} & & 0 \\
& \ddots & \\
0 & & \lambda_{n}
\end{array}\right) \mapsto \lambda_{j}-\lambda_{i} \neq 0 \text { since } i \neq j
$$

So this 1-dimensional representation of $H$ is a weight and $\mathbb{C} E_{i j}$ lies in the weight space.

$$
\operatorname{dim}(H)+\sum \operatorname{dim}(\text { weight space })=\operatorname{dim}(L)=\operatorname{dim}(H)+\sum \operatorname{dim}\left(\mathbb{C} E_{i j}\right)
$$

So $\mathbb{C} E_{i j}$ is the full weight space, giving $L=H \oplus\left(\oplus_{i \neq j} \mathbb{C} E_{i j}\right)$ as a direct sum of Lie algebras.

Summary. $\mathfrak{s l}_{n}(\mathbb{C})=[n \times n$ matrices of trace 0$] . \mathfrak{s l}_{n}(\mathbb{C})$ is simple.

$$
H=\left\{\left.\left(\begin{array}{lll}
\lambda_{1} & & 0 \\
& \ddots & \\
0 & & \lambda_{n}
\end{array}\right) \right\rvert\, \sum_{i=1}^{n} \lambda_{i}=0\right\}
$$

is a Cartan subalgebra. $L=H \oplus\left(\oplus_{i \neq j} \mathbb{C} E_{i j}\right)$. The roots are

$$
\left(\begin{array}{lll}
\lambda_{1} & & 0 \\
& \ddots & \\
0 & & \lambda_{n}
\end{array}\right) \mapsto \lambda_{j}-\lambda_{i} \text { for } i \neq j
$$

## 8. The Cartan Decomposition

Throughout this chapter, let $L$ be a semisimple Lie algebra and $H$ a Cartan subalgebra of $L$.

$$
\begin{gathered}
L=\oplus_{\lambda} L_{\lambda}, \lambda \text { weights }- \text { the Cartan decomposition } \\
H=L_{0} \\
L=H \oplus\left(\oplus_{\alpha \in \Phi} L_{\alpha}\right)
\end{gathered}
$$

Proposition 8.1. If $x \in L_{\lambda}, y \in L_{\mu}, \lambda \neq-\mu$ then $\langle x, y\rangle=0$.
Proof.

$$
\begin{gathered}
\langle x, y\rangle=\operatorname{tr}(\operatorname{ad} x \operatorname{ad} y) \\
L_{V} \operatorname{ad} x \subseteq L_{V+\lambda} \\
L_{V} \text { ad } x \operatorname{ad} y \subseteq L_{V+\lambda+\mu}
\end{gathered}
$$

So if $\lambda+\mu \neq 0$ then $V+\lambda+\mu \neq V$. Choose a basis of $L$ adapted to a Cartan decomposition. ad $x$ ad $y$ is then represented by a block matrix with zero blocks on the diagonal:

$$
\left(\begin{array}{llll}
0 & & & * \\
& 0 & & \\
& & \ddots & \\
* & & & 0
\end{array}\right)
$$

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

Proposition 8.2. If $\alpha \in \Phi$ then $-\alpha \in \Phi$.
Proof. Suppose if possible that $-\alpha \notin \Phi$. Then $L_{-\alpha}=0$ and $L_{\alpha} \neq 0$. Let $x \in L_{\alpha}$. $\langle x, y\rangle=0$ for all $y \in L_{\lambda}$, for all $\lambda$. Hence $\langle x, y\rangle=0$ for all $y \in L$. But $\langle\cdot, \cdot\rangle$ is nondegenerate, so $x=0$. Thus $L_{\alpha}=0$, a contradiction.

Proposition 8.3. The Killing form of $L$ remains non-degenerate on restriction to $H$; i.e. if $x \in H$ and $\langle x, y\rangle=0$ for all $y \in H$ then $x=0$.

Note. We are saying that the Killing form of $L$ restricted to $H$ is non-degenerate, not that the Killing form of $H$ is non-degenerate. In fact, the Killing form of $H$ is degenerate since $H$ is nilpotent.

Proof. Let $x \in H$ satisfy $\langle x, y\rangle=0$ for all $y \in H$. We also have $\langle x, y\rangle=0$ for all $y \in L_{\lambda}$, $\lambda \neq 0$, by 8.1. So $\langle x, y\rangle=0$ for all $y \in L$. So $x=0$.

Proposition 8.4. $[H H]=0$; i.e. $H$ is abelian.
Proof. Suppose $x \in[H H]$ and let $y \in H$.

$$
\begin{gathered}
\langle x, y\rangle=\operatorname{tr}(\operatorname{ad} x \operatorname{ad} y) \\
L=\oplus_{\lambda} L_{\lambda}, \text { each } L_{\lambda} \text { an } H \text {-module. }
\end{gathered}
$$

On $L_{\lambda}$ we have ad $x$ represented by

$$
\left(\begin{array}{ccc}
\lambda(x) & & * \\
& \ddots & \\
0 & & \lambda(x)
\end{array}\right)
$$

and ad $y$ represented by

$$
\left(\begin{array}{ccc}
\lambda(y) & & * \\
& \ddots & \\
0 & & \lambda(y)
\end{array}\right)
$$

Hence, $\operatorname{ad} x \operatorname{ad} y$ is represented by

$$
\left(\begin{array}{ccc}
\lambda(x) \lambda(y) & & * \\
& \ddots & \\
0 & & \lambda(x) \lambda(y)
\end{array}\right)
$$

So $\langle x, y\rangle=\operatorname{tr}(\operatorname{ad} x \operatorname{ad} y)=\sum_{\lambda} \operatorname{dim}\left(L_{\lambda}\right) \lambda(x) \lambda(y)$. Let $h_{1}, h_{2} \in H$. Then

$$
\lambda\left[h_{1} h_{2}\right]=\lambda\left(h_{1}\right) \lambda\left(h_{2}\right)-\lambda\left(h_{2}\right) \lambda\left(h_{1}\right)=0
$$

Since $\lambda$ is a 1 -dimensional representation. Hence $\lambda(x)=0$ for all $x \in[H H]$, so $\langle x, y\rangle=0$ for all $y \in H$. Since $\langle\cdot, \cdot\rangle$ is non-degenerate, $x=0$. So $[H H]=0$; i.e. $H$ is abelian.

Let $H^{*}$ denote the dual space of $H, \operatorname{Hom}(H, \mathbb{C})$. Then $\operatorname{dim}\left(H^{*}\right)=\operatorname{dim}(H)$. We can define a map $H \rightarrow H^{*}, h \mapsto h^{*}$, by $h^{*}(x)=\langle h, x\rangle$ for $x \in H$.

Lemma 8.5. The map $h \mapsto h^{*}$ is an isomorphism of vector spaces.
Proof. The map is clearly linear. Suppose $h$ is in the kernel of this map, i.e. $h^{*}=0$. Then $\langle h, x\rangle=0$ for all $x \in H$. Since the Killing form is non-degenerate, $h=0$. So the kernel is trivial.

$$
\operatorname{dim}(\text { image })=\operatorname{dim}(H)-\operatorname{dim}(\text { kernel })=\operatorname{dim}(H)=\operatorname{dim}\left(H^{*}\right)
$$

Let $\alpha \in \Phi$. Then $\alpha \in H^{*}$, so there is a unique $h_{\alpha} \in H$ such that $\alpha(x)=\left\langle h_{\alpha}, x\right\rangle$ for all $x \in H$.

Proposition 8.6. The $h_{\alpha}$, as defined above, and taken over all $\alpha \in \Phi$, span $H$.
Proof. Suppose not. Then there is an $x \in H, x \neq 0$, such that $\left\langle h_{\alpha}, x\right\rangle=0$ for each $\alpha \in \Phi$. Hence $\alpha(x)=0$ for all $\alpha \in \Phi$. Let $y \in H$.

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

Since $\lambda(x)=0$ for all weights $\lambda,\langle x, y\rangle=0$ for all $y \in H$. So $x=0$, a contradiction.

Proposition 8.7. For $\alpha \in \Phi, h_{\alpha} \in\left[L_{\alpha} L_{-\alpha}\right]$, a subspace of $H$.
Proof. Consider $L_{\alpha}$ as an $H$-module: it contains a 1 -dimensional submodule $\mathbb{C} e_{\alpha}$, $e_{\alpha} \neq 0 .\left[e_{\alpha} x\right]=\alpha(x) e_{\alpha}$ for $x \in H$. Let $y \in L_{-\alpha}$. We show that $\left[y e_{\alpha}\right]=\left\langle y, e_{\alpha}\right\rangle h_{\alpha}$. To see this let $z=\left[y e_{\alpha}\right]-\left\langle y, e_{\alpha}\right\rangle h_{\alpha}$ and let $x \in H$.

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

So $\langle z, x\rangle=0$ for all $x \in H$, so $z=0$. So $\left[y e_{\alpha}\right]=\left\langle y, e_{\alpha}\right\rangle h_{\alpha}$. We can find $y \in L_{-\alpha}$ such that $\left\langle y, e_{\alpha}\right\rangle \neq 0$. If not, $e_{\alpha}$ is orthogonal to $L_{-\alpha}$ and to each $L_{\lambda}$ with $\lambda \neq-\alpha$, so $e_{\alpha} \in L^{\perp}$, so $e_{\alpha}=0$, a contradiction. Choose such a $y$, then

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

Proposition 8.8. Let $\alpha \in \Phi$. Then $\left\langle h_{\alpha}, h_{\alpha}\right\rangle \neq 0$.
Proof. Let $\beta \in \Phi$. We know that $h_{\alpha} \in\left[L_{-\alpha} L_{\alpha}\right]$. There exists $r_{\beta, \alpha} \in \mathbb{Q}$ such that $\beta=r_{\beta, \alpha} \alpha$ on $\left[L_{-\alpha} L_{\alpha}\right]$.

$$
\left\langle h_{\beta}, h_{\alpha}\right\rangle=\beta\left(h_{\alpha}\right)=r_{\beta, \alpha} \alpha\left(h_{\alpha}\right)=r_{\beta, \alpha}\left\langle h_{\alpha}, h_{\alpha}\right\rangle
$$

If $\left\langle h_{\alpha}, h_{\alpha}\right\rangle=0$ then $\left\langle h_{\beta}, h_{\alpha}\right\rangle=0$ for all $\beta \in \Phi$. But the set $\left\{h_{\beta} \mid \beta \in \Phi\right\}$ spans $H$. Hence $\left\langle x, h_{\alpha}\right\rangle=0$ for all $x \in H$. This implies $h_{\alpha}=0$, so $\alpha=0$, a contradiction.

Theorem 8.9. If $\alpha \in \Phi$ then $\operatorname{dim}\left(L_{\alpha}\right)=1$.
Note. Of course, $\operatorname{dim}(H)=\operatorname{dim}\left(L_{0}\right)$ is not generally 1.
Proof. Let $M$ be the subspace of $L$ given by

$$
M=\mathbb{C} e_{\alpha} \oplus \mathbb{C} h_{a} \oplus L_{-\alpha} \oplus L_{-2 \alpha} \oplus \ldots
$$

where $\mathbb{C} e_{\alpha}$ is a 1-dimensional $H$-submodule of $L_{\alpha}$. There are only finitely many summands since $\Phi$ is finite. Recall from the proof of 8.7 that there is an $e_{-\alpha} \in L_{-\alpha}$ such that $\left[e_{-\alpha} e_{\alpha}\right]=h_{\alpha}$. Also, for any $y \in L_{-\alpha},\left[y e_{\alpha}\right]=\left\langle y, e_{\alpha}\right\rangle h_{\alpha}$.

We first show that $M$ ad $e_{\alpha} \subseteq M$ :

$$
\begin{gathered}
{\left[e_{\alpha} e_{\alpha}\right]=0} \\
{\left[h_{\alpha} e_{\alpha}\right]=-\left[e_{\alpha} h_{\alpha}\right]=-\alpha\left(h_{\alpha}\right) e_{\alpha}}
\end{gathered}
$$

Let $y \in L_{-\alpha}$. Then

$$
\begin{gathered}
{\left[y e_{\alpha}\right]=\left\langle y, e_{\alpha}\right\rangle h_{\alpha} \in \mathbb{C} h_{\alpha}} \\
{\left[L_{-i \alpha} e_{\alpha}\right] \subseteq L_{-(i-1) \alpha} \text { for } i \geq 2}
\end{gathered}
$$

Secondly, we show that $M$ ad $e_{-\alpha} \subseteq M$.

$$
\begin{gathered}
{\left[e_{\alpha} e_{-\alpha}\right]=-h_{\alpha}} \\
{\left[h_{\alpha} e_{-\alpha}\right]=-\left[e_{-\alpha} h_{\alpha}\right]=\alpha\left(h_{\alpha}\right) e_{-\alpha} \in L_{-\alpha}} \\
{\left[L_{-i \alpha} e_{-\alpha}\right] \subseteq L_{-(i+1) \alpha} \text { for } i \geq 1}
\end{gathered}
$$

$\left[e_{-\alpha} e_{\alpha}\right]=h_{\alpha}$, so $M \operatorname{ad} h_{a} \subseteq M$. We now calculate $\operatorname{tr}_{M}\left(\operatorname{ad} h_{\alpha}\right)$ in two different ways.

$$
\begin{aligned}
& \operatorname{tr}_{M}\left(\operatorname{ad} h_{\alpha}\right)=\operatorname{tr}_{M}\left(\operatorname{ad} e_{-a} \operatorname{ad} e_{\alpha}-\operatorname{ad} e_{\alpha} \operatorname{ad} e_{-\alpha}\right)=0 \\
& \operatorname{tr}_{M}\left(\operatorname{ad} h_{\alpha}\right)=\alpha\left(h_{\alpha}\right)+0-\operatorname{dim}\left(L_{-\alpha}\right) a\left(h_{\alpha}\right)-2 \operatorname{dim}\left(L_{-2 \alpha}\right) \alpha\left(h_{\alpha}\right)-3 \operatorname{dim}\left(L_{-3 \alpha}\right) \alpha\left(h_{\alpha}\right)-\ldots \\
& =\alpha\left(h_{\alpha}\right)\left(1-\operatorname{dim}\left(L_{-\alpha}\right)-2 \operatorname{dim}\left(L_{-2 \alpha}\right)-\ldots\right) \\
& =\left\langle h_{\alpha}, h_{\alpha}\right\rangle\left(1-\operatorname{dim}\left(L_{-\alpha}\right)-2 \operatorname{dim}\left(L_{-2 \alpha}\right)-\ldots\right)
\end{aligned}
$$

Since $\left\langle h_{\alpha}, h_{\alpha}\right\rangle \neq 0,1-\operatorname{dim}\left(L_{-\alpha}\right)-2 \operatorname{dim}\left(L_{-2 \alpha}\right)-\ldots=0$. Since $-\alpha \in \Phi, \operatorname{dim}\left(L_{-\alpha}\right) \geq 1$. Hence $\operatorname{dim}\left(L_{-\alpha}\right)=1$ and $\operatorname{dim}\left(L_{-i \alpha}\right)=0$ for $i \geq 2$.

So, interchanging $\alpha \leftrightarrow-\alpha$, we have that for each $\alpha \in \Phi, \operatorname{dim}\left(L_{\alpha}\right)=1$ and $\operatorname{dim}\left(L_{i \alpha}\right)=0$ for $i \geq 2, i \in \mathbb{N}$.

We have the following easy corollary:
Corollary 8.10. If $\alpha \in \Phi, m \alpha \in \Phi$ and $m \in \mathbb{Z}$ then $m= \pm 1$.
Now let $\alpha, \beta \in \Phi, \beta \neq \pm \alpha$. Consider

$$
-p \alpha+\beta, \ldots-\alpha+\beta, \beta, \alpha+\beta, 2 \alpha+\beta, \ldots, q \alpha+\beta
$$

There are integers $p, q$ such that all of the above are roots but $-(p+1) \alpha+\beta$ and $(q+1) \alpha+\beta$ are not roots. This collection is called the $\alpha$-chain of roots through $\beta$.

Note that $\beta+(q+1) \alpha, \beta-(p+1) \alpha \neq 0$, so $L_{\beta+(q+1) \alpha}=0, L_{\beta-(p+1) \alpha}=0$. Let

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

Choose $e_{\alpha} \in L_{\alpha}, \quad e_{-\alpha} \in L_{-\alpha}$ such that $\left[e_{-\alpha} e_{\alpha}\right]=h_{\alpha}$. We claim that $M$ ade $e_{\alpha} \subseteq M$, $M \operatorname{ad} e_{-\alpha} \subseteq M$ and $M \operatorname{ad} h_{\alpha} \subseteq M$.

$$
\begin{aligned}
L_{i \alpha+\beta} \text { ad } e_{\alpha} \subseteq L_{(i+1)_{\alpha+\beta}} \\
L_{i \alpha+\beta} \text { ad } e_{-\alpha} \subseteq L_{(i-1)_{\alpha+\beta}}
\end{aligned}
$$

So $M \operatorname{ad} h_{\alpha} \subseteq M$. We now calculate $\operatorname{tr}_{M}\left(\operatorname{ad} h_{\alpha}\right)$ in two different ways:

$$
\begin{gathered}
\operatorname{tr}_{M}\left(\operatorname{ad} h_{\alpha}\right)=\operatorname{tr}_{M}\left(\operatorname{ad} e_{-a} \operatorname{ad} e_{\alpha}-\operatorname{ad} e_{\alpha} \operatorname{ad} e_{-\alpha}\right)=0 \\
\operatorname{tr}_{M}\left(\operatorname{ad} h_{\alpha}\right)=\sum_{i=-p}^{q}(i \alpha+\beta)\left(h_{\alpha}\right) \text { since } \operatorname{dim}\left(L_{i \alpha+\beta}\right)=1
\end{gathered}
$$

So $\sum_{i=-p}^{q}(i \alpha+\beta)\left(h_{\alpha}\right)=0 ;\left(\sum_{i=-p}^{q} i\right) \alpha\left(h_{\alpha}\right)+\sum_{i=-p}^{q} \beta\left(h_{\alpha}\right)=0$.
$\Rightarrow \quad\left(\frac{p(p+1)}{2}-\frac{q(q+1)}{2}\right)\left\langle h_{\alpha}, h_{\alpha}\right\rangle+(p+q+1)\left\langle h_{\beta}, h_{\alpha}\right\rangle=0$
$\Rightarrow \quad\left(\frac{q-p}{2}\right)(p+q+1)\left\langle h_{\alpha}, h_{\alpha}\right\rangle+(p+q+1)\left\langle h_{\beta}, h_{\alpha}\right\rangle=0$
$\Rightarrow \quad 2\left\langle h_{\beta}, h_{\alpha}\right\rangle /\left\langle h_{\alpha}, h_{\alpha}\right\rangle=p-q$
Thus we have proved:
Proposition 8.11. Let $\alpha, \beta \in \Phi, \beta \neq \pm \alpha$. Take the $\alpha$-chain of roots through $\beta$,

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

Then

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

Corollary 8.12. If $\alpha \in \Phi, \xi \alpha \in \Phi$ and $\xi \in \mathbb{C}$ then $\xi= \pm 1$.
Proof. If $\xi \neq \pm 1$ set $\beta=\xi \alpha$. Then

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

So $2 \xi \in \mathbb{Z}$. If $\xi \in \mathbb{Z}$ then $\xi= \pm 1$ by 8.10 . If $\xi \notin \mathbb{Z}$ then $p \not \equiv q$ modulo 2 . The $\alpha$-chain of roots through $\beta$ is

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

But $\beta=\frac{p-q}{2} \alpha$, with $p$ and $q$ not both zero. So $\frac{\alpha}{2}$ appears in the $\alpha$-chain. This implies $\alpha, \frac{\alpha}{2} \in \Phi$, which contradicts 8.10.

Proposition 8.13. For all $\alpha, \beta \in \Phi,\left\langle h_{\alpha}, h_{\beta}\right\rangle \in \mathbb{Q}$.
Proof. If $\beta \neq \pm \alpha$ then by $8.11\left\langle h_{\alpha}, h_{\beta}\right\rangle /\left\langle h_{\alpha}, h_{\alpha}\right\rangle \in \mathbb{Q}$. We show that $\left\langle h_{\alpha}, h_{\alpha}\right\rangle \in \mathbb{Q}$.

$$
\begin{aligned}
\left\langle h_{\alpha}, h_{\alpha}\right\rangle & =\operatorname{tr}\left(\operatorname{ad} h_{\alpha} \operatorname{ad} h_{\alpha}\right) \\
& =\sum_{\beta \in \Phi} \beta\left(h_{\alpha}\right)^{2} \\
& =\sum_{\beta \in \Phi}\left\langle h_{\alpha}, h_{\beta}\right\rangle^{2} \\
\frac{1}{\left\langle h_{\alpha}, h_{\alpha}\right\rangle}= & \sum_{\beta \in \Phi}\left(\frac{\left\langle h_{\alpha}, h_{\beta}\right\rangle}{\left\langle h_{\alpha}, h_{\alpha}\right\rangle}\right)^{2} \in \mathbb{Q}
\end{aligned}
$$

So $\left\langle h_{\alpha}, h_{\alpha}\right\rangle \in \mathbb{Q}$; so $\left\langle h_{\alpha}, h_{\beta}\right\rangle \in \mathbb{Q}$.

## 9. The Root System and the Weyl Group

$\left\{h_{\alpha} \mid \alpha \in \Phi\right\}$ spans $H$, so we can find a subset $\left\{h_{\alpha_{1}}, \ldots, h_{\alpha_{l}}\right\}$ that forms a basis for $H$; $\operatorname{dim}(H)=l$.

Proposition 9.1. Let $\alpha \in \Phi$. Then $h_{\alpha}=\sum_{i=1}^{l} \mu_{i} h_{\alpha_{i}}$ for some $\mu_{i} \in \mathbb{Q}$.

Proof. We know this for $\mu_{i} \in \mathbb{C}$. Let $\left\langle h_{\alpha_{i}}, h_{\alpha_{j}}\right\rangle=\xi_{i j} \in \mathbb{Q}$. The matrix $\Xi=\left(\xi_{i j}\right)$ is nonsingular, since if it were singular we would have $\eta_{1}, \ldots, \eta_{l}$ not all zero such that $\sum_{i=1}^{l} \eta_{i} \xi_{i j}=0$. Then

$$
\left\langle\sum_{i=1}^{l} \eta_{i} h_{\alpha_{i}}, h_{\alpha_{i}}\right\rangle=\sum_{i=1}^{l} \eta_{i} \xi_{i j}=0
$$

So $\left\langle\sum_{i=1}^{l} \eta_{i} h_{\alpha_{i}}, x\right\rangle=0$ for all $x \in H .\langle\cdot \cdot\rangle$ is non-degenerate, so all $\eta_{i}=0$, which is a contradiction.

$$
\begin{gathered}
\left\langle h_{\alpha}, h_{\alpha_{1}}\right\rangle=\mu_{1} \xi_{11}+\ldots+\mu_{l} \xi_{l 1} \\
\vdots \\
\left\langle h_{\alpha}, h_{\alpha_{l}}\right\rangle=\mu_{1} \xi_{1 l}+\ldots+\mu_{l} \xi_{l l}
\end{gathered}
$$

We have $l$ linear equations in $l$ unknowns with a non-singular coefficient matrix, all the entries of which are rational. Hence, by Cramer's Rule, there is a unique solution $\left(\mu_{i}\right) \in \mathbb{Q}^{l}$

Let $H_{\mathbb{Q}}$ be the set of all $\sum_{i=1}^{l} \mu_{i} h_{\alpha_{i}}, \mu_{i} \in \mathbb{Q} . \operatorname{dim}_{\mathbb{Q}}\left(H_{\mathbb{Q}}\right)=l . H_{\mathbb{Q}}$ is independent of the choice of basis; all $h_{\alpha} \in H_{\mathbb{Q}}$.

Let $H_{\mathbb{R}}$ be the set of all $\sum_{i=1}^{l} \mu_{i} h_{\alpha_{i}}, \mu_{i} \in \mathbb{R} . \operatorname{dim}_{\mathbb{R}}\left(H_{\mathbb{R}}\right)=l$.

Proposition 9.2. Let $x \in H_{\mathbb{R}}$. Then $\langle x, x\rangle \in \mathbb{R}_{\geq 0}$ and $\langle x, x\rangle=0 \Leftrightarrow x=0$.

Proof. Let $x \in H_{\mathbb{R}}, x=\sum_{i=1}^{l} \mu_{i} h_{\alpha_{i}}$.

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

So $\langle x, x\rangle \in \mathbb{R}$ and $\langle x, x\rangle \geq 0$. Suppose $\langle x, x\rangle=0$. Then for all $\alpha \in \Phi, \sum_{i} \mu_{i} \alpha\left(h_{\alpha_{i}}\right)=0$. In particular, $\sum_{i} \mu_{i} \alpha_{j}\left(h_{\alpha_{i}}\right)=0$ for $j=1, \ldots, l ; \sum_{i} \mu_{i}\left\langle h_{\alpha_{i}}, h_{\alpha_{j}}\right\rangle=\sum_{i} \mu_{i} \xi_{i j}=0$ for all $j$. $\Xi$ is non-singular, so $\mu_{i}=0$ for all $i$, so $x=0$.

So all $h_{\alpha} \in H_{\mathbb{R}} ; \operatorname{dim}_{\mathbb{R}}\left(H_{\mathbb{R}}\right)=l$. We introduce a total order on $H_{\mathbb{R}}$ : let $x \in H_{\mathbb{R}}$, $x=\sum_{i} \mu_{i} h_{\alpha_{i}}$. If $x \neq 0$ we say $x \succ 0$ if the first non-zero $\mu_{i}$ is positive; if $x \neq 0$ we say $x \prec 0$ if the first non-zero $\mu_{i}$ is negative. We have trichotomy: for each $x \in H_{\mathbb{R}}$ precisely one of $x=0, x \prec 0, x \succ 0$ is true.

So, for $\alpha \in \Phi, h_{\alpha} \prec 0$ or $h_{\alpha} \succ 0$. Define $\alpha \prec 0$ if $h_{\alpha} \prec 0$ and $\alpha \succ 0$ if $h_{\alpha} \succ 0$. Define

$$
\begin{gathered}
\Phi^{+}=\{\alpha \in \Phi \mid \alpha \succ 0\}, \text { the positive roots, and } \\
\Phi^{-}=\{\alpha \in \Phi \mid \alpha \prec 0\}, \text { the negative roots. }
\end{gathered}
$$

Clearly, $\Phi=\Phi^{+} \cup \Phi^{-}$.
A fundamental root is a positive root that is not the sum of two positive roots. Let $\Pi$ be the set of fundamental roots.

Proposition 9.3. (i) Every positive root is a sum of fundamental roots.
(ii) $\left\{h_{\alpha} \mid \alpha \in \Pi\right\}$ is a basis of $H_{\mathbb{R}}$.
(iii) If $\alpha, \beta \in \Pi$ and $\alpha \neq \beta$ then $\left\langle h_{\alpha}, h_{\beta}\right\rangle \leq 0$.

Proof. (i) Let $\alpha \in \Phi^{+}$. If $\alpha \in \Pi$ we are done. If $\alpha \notin \Pi$ the there exist $\beta, \gamma \in \Phi^{+}$such that $\alpha=\beta+\gamma$ with $\beta, \gamma \prec \alpha$. Repeat to get the result.
(iii) Let $\alpha, \beta \in \Pi$ with $\alpha \neq \beta$. Then $\alpha-\beta \notin \Phi$ since if not

$$
\alpha=(\alpha-\beta)+\beta \text { or } \beta=(\beta-\alpha)+\alpha
$$

so either $\alpha$ or $\beta$ would be a sum of positive roots. Consider the $\alpha$-chain of roots through $\beta$ :

$$
\begin{aligned}
& \beta, \alpha+\beta, \ldots, q \alpha+\beta \\
\Rightarrow \quad & 2 \frac{\left\langle h_{\alpha}, h_{\beta}\right\rangle}{\left\langle h_{\alpha}, h_{\alpha}\right\rangle}=p-q=-q
\end{aligned}
$$

$$
\Rightarrow \quad\left\langle h_{\alpha}, h_{\beta}\right\rangle \leq 0 \quad \text { since }\left\langle h_{\alpha}, h_{\alpha}\right\rangle \geq 0 \text { by 9.2. }
$$

(ii) By (i), the $h_{\alpha}$ for $\alpha \in \Phi$ span $H$. We show the $h_{\alpha}$ are linearly independent. Suppose not: then there exist $\mu_{i} \in \mathbb{R}$ not all zero such that

$$
\sum_{\alpha_{i} \in \Pi} \mu_{i} h_{\alpha_{i}}=0
$$

Rearrange this sum, taking all the positive $\mu_{i}$ to one side. Then

$$
\begin{gathered}
x=\mu_{i_{1}} h_{\alpha_{i_{1}}}+\ldots+\mu_{i_{r}} h_{\alpha_{i_{r}}}=\mu_{j_{1}} h_{\alpha_{j_{1}}}+\ldots+\mu_{j_{s}} h_{\alpha_{s_{s}}} \\
\mu_{i_{u}}, \mu_{j_{v}}>0, i_{u}, j_{v} \text { distinct for } 1 \leq u \leq r, 1 \leq v \leq s .
\end{gathered}
$$

Then

$$
\langle x, x\rangle=\left\langle\mu_{i_{1}} h_{\alpha_{i_{1}}}+\ldots+\mu_{i_{r}} h_{\alpha_{i_{r}}}, \mu_{j_{1}} h_{\alpha_{j_{1}}}+\ldots+\mu_{j_{s}} h_{\alpha_{j_{s}}}\right\rangle \leq 0
$$

by (iii). So $x=0$, a contradiction.

Note. $\Phi^{+}$can be chosen in many different ways. However, $\Pi$ is determined by $\Phi^{+}$and $\Phi^{+}$is determined by $\Pi$.

Example. Let $L=\mathfrak{s l}_{n}(\mathbb{C})$. The roots are

$$
\left(\begin{array}{lll}
\lambda_{1} & & 0 \\
& \ddots & \\
0 & & \lambda_{n}
\end{array}\right) \mapsto \lambda_{j}-\lambda_{i} \text { for } j \neq i
$$

Define $\Phi^{+}$to be the roots with $j>i$. Then the fundamental roots are

$$
\begin{gathered}
\left(\begin{array}{ccc}
\lambda_{1} & & 0 \\
& \ddots & \\
0 & & \lambda_{n}
\end{array}\right) \mapsto \lambda_{i+1}-\lambda_{i} \text { for } 1 \leq i \leq n-1 \\
\lambda_{j}-\lambda_{i}=\left(\lambda_{i+1}-\lambda_{i}\right)+\left(\lambda_{i+2}-\lambda_{i+1}\right)+\ldots+\left(\lambda_{j}-\lambda_{j-1}\right) \\
\operatorname{dim}(H)=n-1=l, \text { the rank of } L .
\end{gathered}
$$

For each $\alpha \in \Phi$ we define $s_{\alpha}: H_{\mathbb{R}} \rightarrow H_{\mathbb{R}}$ by

$$
s_{\alpha}(x)=x-2 \frac{\left\langle h_{\alpha}, x\right\rangle}{\left\langle h_{\alpha}, h_{\alpha}\right\rangle} h_{\alpha}
$$

$s_{\alpha}$ is linear and $s_{\alpha}\left(h_{\alpha}\right)=-h_{\alpha}$. The set of $x$ such that $\langle x, x\rangle=0$ forms a hyperplane i.e. a subspace of codimension 1. $s_{\alpha}$ is the reflection of $H_{\mathbb{R}}$ in the hyperplane orthogonal to $h_{\alpha}$.

$$
\begin{gathered}
s_{\alpha}^{2}=\mathrm{id} \\
s_{\alpha}=s_{-\alpha}
\end{gathered}
$$

Let $W$ be the group of all non-singular linear maps $H_{\mathbb{R}} \rightarrow H_{\mathbb{R}}$ generated by $\left\{s_{\alpha} \mid \alpha \in \Phi\right\} . W$ is called the Weyl group. ${ }^{\dagger}$

Proposition 9.4. (i) $W$ is a finite group.
(ii) $W$ is a group of isometries, i.e. for all $x, y \in H_{\mathbb{R}}, w \in W,\langle w(x), w(y)\rangle=\langle x, y\rangle$.
(iii) For each $\alpha \in \Phi$ and $w \in W$ there is a $\beta \in \Phi$ such that $w\left(h_{\alpha}\right)=h_{\beta}$.

Proof. (ii) Let $x, y \in H_{\mathbb{R}}$. Then

$$
\begin{aligned}
\left\langle s_{\alpha}(x), s_{\alpha}(y)\right\rangle & =\left\langle x-2 \frac{\left\langle h_{\alpha}, x\right\rangle}{\left\langle h_{\alpha}, h_{\alpha}\right\rangle} h_{\alpha}, y-2 \frac{\left\langle h_{\alpha}, y\right\rangle}{\left\langle h_{\alpha}, h_{\alpha}\right\rangle} h_{\alpha}\right\rangle \\
& =\langle x, y\rangle-4 \frac{\left\langle h_{\alpha}, x\right\rangle\left\langle h_{\alpha}, y\right\rangle}{\left\langle h_{\alpha}, h_{\alpha}\right\rangle}-4 \frac{\left\langle h_{\alpha}, x\right\rangle\left\langle h_{\alpha}, y\right\rangle}{\left\langle h_{\alpha}, h_{\alpha}\right\rangle^{2}}\left\langle h_{\alpha}, h_{\alpha}\right\rangle \\
& =\langle x, y\rangle
\end{aligned}
$$

So $s_{\alpha}$ is an isometry; so $w$ is an isometry for all $w \in W$.
(iii) Now consider $s_{\alpha}\left(h_{\beta}\right)$ :

[^2]\[

$$
\begin{aligned}
& s_{\alpha}\left(h_{\alpha}\right)=h_{-\alpha} \\
& s_{\alpha}\left(h_{-\alpha}\right)=h_{\alpha}
\end{aligned}
$$
\]

So suppose $\beta \neq \pm \alpha$. Consider the $\alpha$-chain of roots through $\beta$,

$$
\begin{aligned}
-p \alpha & +\beta, \ldots, \beta, \ldots, q \alpha+\beta \\
s_{\alpha}\left(h_{\beta}\right) & =h_{\beta}-2 \frac{\left\langle h_{\alpha}, h_{\beta}\right\rangle}{\left\langle h_{\alpha}, h_{\alpha}\right\rangle} h_{\alpha} \\
& =h_{\beta}-(p-q) h_{\alpha} \\
& =h_{\beta+(q-p) \alpha}
\end{aligned}
$$

Now $\beta+(q-p) \alpha \in \Phi$ since $-p \leq p-q \leq q$. So $s_{\alpha}$ permutes the $h_{\beta}$ for $\beta \in \Phi$. Hence $w \in W$ permutes the $h_{\beta}$ for $\beta \in \Phi$. Note that

$$
\beta+((q-p) \alpha+\beta)=(-p \alpha+\beta)+(q \alpha+\beta)
$$

so $s_{\alpha}$ inverts the $h_{\beta}$ in a given $\alpha$-chain.
(i) We have a homomorphism from $W$ to the group of permutations of the $h_{\alpha}$ for $\alpha \in \Phi . \Phi$ is finite, so the image of this homomorphism is finite. If $w \in W$ is in the kernel then $w\left(h_{\alpha}\right)=h_{\alpha}$ for all $\alpha \in \Phi$. Since the $h_{\alpha}$ span $H_{\mathbb{R}}, w=\mathrm{id}$. Hence, $W$ is finite.

Proposition 9.5. Given any root $\alpha \in \Phi$ there exists a fundamental root $\alpha_{i} \in \Pi$ and $a$ $w \in W$ such that $h_{\alpha}=w\left(h_{\alpha_{i}}\right)$.

Proof. Each $\alpha \in \Phi$ has the form $\alpha=n_{1} \alpha_{1}+\ldots+n_{l} \alpha_{l}, n_{i} \in \mathbb{Z}$. If $\alpha \in \Phi^{+}$then all $n_{i} \geq 0$; if $\alpha \in \Phi^{-}$then all $n_{i} \leq 0$. We may assume $\alpha \in \Phi^{+}$since if $\alpha \in \Phi^{-}$then use $h_{\alpha}=s_{\alpha}\left(h_{-\alpha}\right)$. The quantity $n_{1}+\ldots+n_{l}$ is called the height of $\alpha, \operatorname{ht}(\alpha)$. We use induction on $\operatorname{ht}(\alpha)$. If $\operatorname{ht}(\alpha)=1$ we are done, so assume $\operatorname{ht}(\alpha)>1$. By 8.12 , at least two $n_{i}>0$.

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

All $n_{i} \geq 0$, so there exists $i$ such that $\left\langle h_{\alpha_{i}}, h_{\alpha}\right\rangle>0$. Let $s_{\alpha_{i}}\left(h_{\alpha}\right)=h_{\beta}$.

$$
\begin{aligned}
h_{\beta} & =h_{\alpha}-2 \frac{\left\langle h_{\alpha_{i}}, h_{\alpha}\right\rangle}{\left\langle h_{\alpha_{i}}, h_{\alpha_{i}}\right\rangle} h_{\alpha_{i}} \\
\beta & =\alpha-2 \frac{\left\langle h_{\alpha_{i}}, h_{\alpha}\right\rangle}{\left\langle h_{\alpha_{i}}, h_{\alpha_{i}}\right\rangle} \alpha_{i}
\end{aligned}
$$

So $\operatorname{ht}(\beta)<\operatorname{ht}(\alpha)$. Passing from $\alpha$ to $\beta$ changes only one $n_{i}$, hence $\beta$ has at least one $n_{j}>0$, so $\beta \in \Phi^{+}$. By induction, $\beta=w^{\prime}\left(h_{\alpha_{j}}\right)$ for some $w^{\prime} \in W$ and some $\alpha_{j} \in \Pi$. Thus, taking $w=s_{\alpha_{i}} w^{\prime} \in W$,

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

Proposition 9.6. The Weyl group $W$ is generated by $s_{\alpha_{1}}, \ldots, s_{\alpha_{l}}$ for $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$.
Proof. Suppose $W_{0}$ is the subgroup generated by $\left\{s_{\alpha_{i}} \mid \alpha_{i} \in \Pi\right\}$. To show $W=W_{0}$ we show $s_{\alpha} \in W_{0}$ for all $\alpha \in \Phi$. The proof of 9.5 shows that $h_{\alpha}=w\left(h_{\alpha_{i}}\right)$ for some $\alpha_{i} \in \Pi$ and some $w \in W$. Consider $w s_{\alpha_{i}} w^{-1} \in W_{0}$.

$$
w s_{\alpha_{i}} w^{-1}\left(h_{\alpha}\right)=w s_{\alpha_{i}}\left(h_{\alpha_{i}}\right)=w\left(-h_{\alpha_{i}}\right)=-h_{\alpha}
$$

Let $x \in H_{\mathbb{R}}$ be such that $\left\langle h_{\alpha}, x\right\rangle=0$. Then

$$
\begin{array}{lc}
\Rightarrow & \left\langle w^{-1}\left(h_{\alpha}\right), w^{-1}(x)\right\rangle=0 \\
\Rightarrow & \left\langle h_{\alpha_{i}}, w^{-1}(x)\right\rangle=0 \\
\Rightarrow & w s_{\alpha_{i}} w^{-1}(x)=w w^{-1}(x)=x
\end{array}
$$

Hence, $w s_{\alpha_{i}} w^{-1}=s_{\alpha}$. Then $s_{\alpha} \in W_{0}$, so $W=W_{0}$.

Example. $L=\mathfrak{s l}_{3}(\mathbb{C}) ; \operatorname{dim}(L)=8$.

$$
H=\left\{\left.\left(\begin{array}{ccc}
\lambda_{1} & 0 & 0 \\
0 & \lambda_{2} & 0 \\
0 & 0 & \lambda_{3}
\end{array}\right) \right\rvert\, \lambda_{1}+\lambda_{2}+\lambda_{3}=0\right\}
$$

$\operatorname{dim}(H)=2$.

$$
L=H \oplus \mathbb{C} E_{12} \oplus \mathbb{C} E_{23} \oplus \mathbb{C} E_{13} \oplus \mathbb{C} E_{21} \oplus \mathbb{C} E_{32} \oplus \mathbb{C} E_{31}
$$

Let

$$
\begin{gathered}
h=\left(\begin{array}{ccc}
\lambda_{1} & 0 & 0 \\
0 & \lambda_{2} & 0 \\
0 & 0 & \lambda_{3}
\end{array}\right) \\
{\left[E_{i j} h\right]=\left(\lambda_{j}-\lambda_{i}\right) E_{i j}}
\end{gathered}
$$

The roots are

$$
\begin{array}{ccc}
\alpha_{1}: h \mapsto \lambda_{2}-\lambda_{1} & \alpha_{2}: h \mapsto \lambda_{3}-\lambda_{2} & \alpha_{1}+\alpha_{2}: h \mapsto \lambda_{3}-\lambda_{1} \\
-\alpha_{1}: h \mapsto \lambda_{1}-\lambda_{2} & -\alpha_{2}: h \mapsto \lambda_{2}-\lambda_{3} & -\alpha_{1}-\alpha_{2}: h \mapsto \lambda_{1}-\lambda_{3} \\
& \Phi=\left\{ \pm \alpha_{1}, \pm \alpha_{2}, \pm\left(\alpha_{1}+\alpha_{2}\right)\right\} & \\
\Pi=\left\{\alpha_{1}, \alpha_{2}\right\} &
\end{array}
$$

Consider the corresponding vectors $h_{\alpha} \in H$. Let

$$
\begin{aligned}
& \quad h=\left(\begin{array}{ccc}
\lambda_{1} & 0 & 0 \\
0 & \lambda_{2} & 0 \\
0 & 0 & \lambda_{3}
\end{array}\right) \in H, h^{\prime}=\left(\begin{array}{ccc}
\mu_{1} & 0 & 0 \\
0 & \mu_{2} & 0 \\
0 & 0 & \mu_{3}
\end{array}\right) \in H \\
& \left\langle h, h^{\prime}\right\rangle=\operatorname{tr}\left(\operatorname{ad} h \operatorname{ad} h^{\prime}\right) \\
& = \\
& =2\left(\lambda_{2}-\lambda_{1}\right)\left(\mu_{2}-\mu_{1}\right)+2\left(\lambda_{3}-\lambda_{2}\right)\left(\mu_{3}-\mu_{2}\right)+2\left(\lambda_{3}-\lambda_{1}\right)\left(\mu_{3}-\mu_{1}\right) \\
& =2\left(2 \lambda_{1} \mu_{1}+2 \lambda_{2} \mu_{2}+2 \lambda_{3} \mu_{3}-\left(\lambda_{1} \mu_{2}+\lambda_{1} \mu_{3}+\lambda_{2} \mu_{1}+\lambda_{2} \mu_{3}+\lambda_{3} \mu_{1}+\lambda_{3} \mu_{2}\right)\right) \\
& =4\left(\lambda_{1} \mu_{1}+\lambda_{2} \mu_{2}+\lambda_{3} \mu_{3}\right)-2\left(\lambda_{1}+\lambda_{2}+\lambda_{3}\right)\left(\mu_{1}+\mu_{2}+\mu_{3}\right)+2\left(\lambda_{1} \mu_{1}+\lambda_{2} \mu_{2}+\lambda_{3} \mu_{3}\right) \\
& = \\
& =6\left(\lambda_{1} \mu_{1}+\lambda_{2} \mu_{2}+\lambda_{3} \mu_{3}\right) \\
&
\end{aligned}
$$

$h_{\alpha_{1}}$ satisfies $\left\langle h_{\alpha_{1}}, h\right\rangle=\alpha_{1}(h)=\lambda_{2}-\lambda_{1}$, so

$$
h_{\alpha_{1}}=\frac{1}{6}\left(\begin{array}{ccc}
-1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & 0
\end{array}\right)
$$

Similarly,

$$
h_{\alpha_{2}}=\frac{1}{6}\left(\begin{array}{ccc}
0 & 0 & 0 \\
0 & -1 & 0 \\
0 & 0 & 1
\end{array}\right)
$$

For $x \in H_{\mathbb{R}}$ define $|x|=\sqrt{\langle x, x\rangle}$. With this notation, $\left|h_{\alpha_{1}}\right|=\left|h_{\alpha_{2}}\right|=1 / \sqrt{3}$ and

$$
\left\langle h_{\alpha_{1}}, h_{\alpha_{2}}\right\rangle=6 \frac{1}{6} \frac{1}{6}(-1)=-\frac{1}{6}
$$

The angle between $h_{\alpha_{1}}$ and $h_{\alpha_{2}}$ is given by the cosine formula:

$$
\begin{aligned}
\left\langle h_{\alpha_{1}}, h_{\alpha_{2}}\right\rangle & =\left|h_{\alpha_{1}} \| h_{\alpha_{2}}\right| \cos \theta \\
\theta & =2 \pi / 3
\end{aligned}
$$



$$
W=\left\{\mathrm{id}, s_{\alpha_{1}}, s_{\alpha_{2}}, s_{\alpha_{1}} s_{\alpha_{2}}, s_{\alpha_{2}} s_{\alpha_{1}}, s_{\alpha_{1}} s_{\alpha_{2}} s_{\alpha_{1}}=s_{\alpha_{2}} s_{\alpha_{1}} s_{\alpha_{2}}=s_{\alpha_{1}+\alpha_{2}}\right\}
$$

## 10. The Dynkin Diagram

We shall consider the geometrical properties of the $h_{\alpha}$ for $\alpha \in \Phi$.
Proposition 10.1. Let $\alpha, \beta \in \Phi, \beta \neq \pm \alpha$. Then
(i) the angle between $\alpha$ and $\beta$ is one of

$$
\pi / 6, \pi / 4, \pi / 3, \pi / 2,2 \pi / 3,3 \pi / 4,5 \pi / 6
$$

(ii) if the angle is $\pi / 3$ or $2 \pi / 3, h_{\alpha}$ and $h_{\beta}$ have the same length;
(iii) if the angle is $\pi / 4$ or $3 \pi / 4$, the ratio of the lengths of $h_{\alpha}$ and $h_{\beta}$ is $\sqrt{2}$;
(iv) if the angle is $\pi / 6$ or $5 \pi / 6$, the ratio of the lengths of $h_{\alpha}$ and $h_{\beta}$ is $\sqrt{3}$.

Proof. Let $\theta_{\alpha \beta}$ be the angle between $h_{\alpha}$ and $h_{\beta}$. We have

$$
\begin{array}{cc} 
& \left\langle h_{\alpha}, h_{\beta}\right\rangle=\left|h_{\alpha}\right|\left|h_{\beta}\right| \cos \theta_{\alpha \beta} \\
\Rightarrow & \left\langle h_{\alpha}, h_{\beta}\right\rangle^{2}=\left\langle h_{\alpha}, h_{\alpha}\right\rangle\left\langle h_{\beta}, h_{\beta}\right\rangle \cos ^{2} \theta_{\alpha \beta} \\
\Rightarrow \quad 4 \cos ^{2} \theta_{\alpha \beta}=\frac{2\left\langle h_{\alpha}, h_{\beta}\right\rangle}{\left\langle h_{\alpha}, h_{\alpha}\right\rangle} \frac{2\left\langle h_{\beta}, h_{\alpha}\right\rangle}{\left\langle h_{\beta}, h_{\beta}\right\rangle}
\end{array}
$$

By 8.11 , both factors on the RHS are integers, so $4 \cos ^{2} \theta_{\alpha \beta} \in \mathbb{Z} .0 \leq \cos ^{2} \theta_{\alpha \beta}<1$, so $0 \leq 4 \cos ^{2} \theta_{\alpha \beta}<4$, so $4 \cos ^{2} \theta_{\alpha \beta} \in\{0,1,2,3\}$.

$$
\begin{array}{cc}
\Rightarrow & \cos ^{2} \theta_{\alpha \beta} \in\left\{0, \pm \frac{1}{2}, \pm \frac{1}{\sqrt{2}}, \pm \frac{\sqrt{3}}{2}\right\} \\
\Rightarrow & \theta_{\alpha \beta} \in\{\pi / 2, \pi / 3,2 \pi / 3, \pi / 4,3 \pi / 4, \pi / 6,5 \pi / 6\} \\
& 4 \cos ^{2} \theta_{\alpha \beta}=\left(2 \frac{\left\langle h_{\alpha}, h_{\beta}\right\rangle}{\left\langle h_{\alpha}, h_{\alpha}\right\rangle}\right)\left(2 \frac{\left\langle h_{\beta}, h_{\alpha}\right\rangle}{\left\langle h_{\beta}, h_{\beta}\right\rangle}\right)
\end{array}
$$

Suppose $\theta_{\alpha \beta}$ is $\pi / 3$ or $2 \pi / 3$, so $4 \cos ^{2} \theta_{\alpha \beta}=1.1=1 \cdot 1=(-1) .(-1)$, so $\left|h_{\alpha}\right|=\left|h_{\beta}\right|$.
Suppose $\theta_{\alpha \beta}$ is $\pi / 4$ or $3 \pi / 4$, so $4 \cos ^{2} \theta_{\alpha \beta}=2 \cdot 2=1.2=2 \cdot 1=(-1) \cdot(-2)=(-2) \cdot(-1)$. So one of $\left\langle h_{\alpha}, h_{\alpha}\right\rangle$ and $\left\langle h_{\beta}, h_{\beta}\right\rangle$ is twice the other, so one of $\left|h_{\alpha}\right|,\left|h_{\beta}\right|$ is $\sqrt{2}$ times the other.

Suppose $\theta_{\alpha \beta}$ is $\pi / 6$ or $5 \pi / 6$, so $4 \cos ^{2} \theta_{\alpha \beta}=3 \cdot 3=1 \cdot 3=3 \cdot 1=(-1) \cdot(-3)=(-3) \cdot(-1)$. So, as above, one of $\left|h_{\alpha}\right|,\left|h_{\beta}\right|$ is $\sqrt{3}$ times the other.

Proposition 10.2. Let $\alpha \in \Phi$. Then every $\alpha$-chain of roots has at most four roots in it.
Proof. Consider the $\alpha$ - chain of roots through $\beta$ with $\beta$ as the first root:

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

By $8.11,2\left\langle h_{\alpha}, h_{\beta}\right\rangle /\left\langle h_{\alpha}, h_{\alpha}\right\rangle=-q$. The LHS is $0,-1,-2$ or -3 by 10.1 . So $q \leq 3$. So the length of the $\alpha$-chain is at most 4 .

Let

$$
a_{i j}=2 \frac{\left\langle h_{\alpha_{i}}, h_{\alpha_{j}}\right\rangle}{\left\langle h_{\alpha_{i}}, h_{\alpha_{i}}\right\rangle}
$$

and $A=\left(a_{i j}\right) . A$ is called the Cartan matrix; the $a_{i j}$ are the Cartan integers.
Proposition 10.3. The Cartan matrix has the following properties:
(i) for each $i, a_{i i}=2$;
(ii) for $i \neq j, a_{i j} \in\{0,-1,-2,-3\}$;
(iii) $a_{i j}=-2 \Rightarrow a_{j i}=-1 ; a_{i j}=-3 \Rightarrow a_{j i}=-1$;
(iv) $a_{i j}=0 \Leftrightarrow a_{j i}=0$.

Proof. If $i \neq j$ then $4 \cos ^{2} \theta_{\alpha \beta}=a_{i j} a_{j i}$.
(i) Clear.
(ii) Follows from 10.1, 9.3.
(iii) Follows from 10.1.
(iv) Clear.

We incorporate this information into a graph. The Dynkin ${ }^{\dagger}$ diagram is a graph $\Delta$ with $l$ vertices, one for each fundamental root. If $i \neq j$ then vertices $i, j$ are joined by $n_{i j}=a_{i j} a_{j i}$ edges, $0 \leq n_{i j} \leq 3$. The Dynkin diagram may be disconnected, as in


[^3]It splits into connected components and the Cartan matrix splits into corresponding blocks; off-diagonal blocks are zero:

$$
A=\left(\begin{array}{c|c|c}
* & & 0 \\
\hline & \ddots & \\
\hline 0 & & *
\end{array}\right)
$$

We define a corresponding quadratic form $Q$ :

$$
Q\left(x_{1}, \ldots, x_{l}\right)=\sum_{i=1}^{l} 2 x_{i}^{2}-\sum_{1 \leq i \neq j \leq l} \sqrt{n_{i j}} x_{i} x_{j}
$$

Recall the correspondence between quadratic forms on $\mathbb{R}$ and real symmetric matrices:

$$
\begin{gathered}
M=\left(m_{i j}\right) \text { symmetric } \\
x M x^{\mathrm{T}}=\sum_{i, j} m_{i j} x_{i} x_{j}
\end{gathered}
$$

The matrix of $Q\left(x_{1}, \ldots, x_{l}\right)$ is

$$
\left(\begin{array}{ccccc}
2 & -\sqrt{n_{12}} & -\sqrt{n_{13}} & & \\
-\sqrt{n_{12}} & 2 & -\sqrt{n_{23}} & \ddots & \\
-\sqrt{n_{13}} & -\sqrt{n_{23}} & 2 & \ddots & \\
& \ddots & \ddots & \ddots & \\
& & & & 2
\end{array}\right)
$$

Proposition 10.4. The quadratic form $Q\left(x_{1}, \ldots, x_{l}\right)$ is positive definite, i.e. $Q\left(x_{1}, \ldots, x_{l}\right) \geq 0$ and $Q\left(x_{1}, \ldots, x_{l}\right)=0 \Leftrightarrow x_{1}=\ldots=x_{l}=0$.

## Proof.

$$
\begin{gathered}
4 \cos ^{2} \theta_{i j}=a_{i j} a_{j i}=n_{i j} \\
2 \cos \theta_{i j}=-\sqrt{n_{i j}} \\
\left\langle h_{\alpha_{i}}, h_{\alpha_{j}}\right\rangle=\left|h_{\alpha_{i}}\right|\left|h_{\alpha_{j}}\right| \cos \theta_{i j} \\
Q\left(x_{1}, \ldots, x_{l}\right)=\sum_{i, j=1}^{l} 2 \frac{\left\langle h_{\alpha_{i}}, h_{\alpha_{j}}\right\rangle}{\left|h_{\alpha_{i}}\right|\left|h_{\alpha_{j}}\right|} x_{i} x_{j}=2\left\langle\sum_{i} \frac{x_{i}}{\left|h_{\alpha_{i}}\right|} h_{\alpha_{i}}, \sum_{j} \frac{x_{j}}{\left|h_{\alpha_{j}}\right|} h_{\alpha_{j}}\right\rangle=2\langle y, y\rangle
\end{gathered}
$$

where $y=\sum_{i} x_{i} h_{\alpha_{i}} /\left|h_{\alpha_{i}}\right|$. So $Q\left(x_{1}, \ldots, x_{l}\right) \geq 0$. If $Q\left(x_{1}, \ldots, x_{l}\right)=0$ then $\langle y, y\rangle=0$, so $y=0$, so all $x_{i}=0$. The converse is clear.

Recall. Any quadratic form can be diagonalized; there exists a non-singular real $l \times l$ matrix $P$ such that $P M P^{\mathrm{T}}=D$, a diagonal matrix. Let $y=x P^{-1}$; then $x M x^{\mathrm{T}}=y D y^{\mathrm{T}}$.

Proposition 10.5. Let $M=\left(m_{i j}\right)$ be an $l \times l$ real symmetric matrix. Then the associated quadratic form $\sum_{i, j} m_{i j} x_{i} x_{j}$ is positive definite if and only if all leading minors of $M$ have positive determinant. (The leading minors are

$$
\left.\left(m_{11}\right),\left(\begin{array}{ll}
m_{11} & m_{12} \\
m_{21} & m_{22}
\end{array}\right), \ldots, M .\right)
$$

Proof. We use induction on $l$. Assume the quadratic form is positive definite. If $l=1$, $M=\left(m_{11}\right) . m_{11} x^{2}>0 \Leftrightarrow m_{11}>0$. Suppose $l>1 . \sum_{i, j=1}^{l} m_{i j} x_{i} x_{j}$ is still positive definite as it is the original with $x_{l}=0$. By induction, the first $l-1$ leading minors of $M$ have positive determinant; we require that $\operatorname{det}(M)>0 . x M x^{\mathrm{T}}=y D y^{\mathrm{T}}, D$ diagonal with entries $d_{1}, \ldots, d_{l}>0$. Now if $P M P^{\mathrm{T}}=D$,

$$
\operatorname{det}(P)^{2} \operatorname{det}(M)=\operatorname{det}(D)>0 .
$$

Conversely, suppose that all leading minors of $M$ have positive determinant. The same is true of the smaller $(l-1) \times(l-1)$ leading minor. By induction, $\sum_{i, j=1}^{l-1} m_{i j} x_{i} x_{j}$ is positive definite. So we have a diagonal form in new coordinates $y_{1}, \ldots, y_{l}$ :

$$
\begin{gathered}
\sum_{i, j=1}^{l-1} m_{i j} x_{i} x_{j}=\sum_{k=1}^{l-1} d_{k} x_{k}^{2} \text { with } d_{k}>0 . \\
\sum_{i, j=1}^{l} m_{i j} x_{i} x_{j}=\sum_{k=1}^{l-1} d_{k} x_{k}^{2}+2 e_{1} y_{1} x_{l}+\ldots+2 e_{l-1} y_{l-1} x_{l}+e x_{l}^{2}
\end{gathered}
$$

This may be diagonalized by a further transformation of coordinates:

$$
z_{i}=y_{i}+\frac{e_{i}}{d_{i}} x_{i}
$$

We get $d_{1} z_{1}^{2}+\ldots d_{l-1} z_{l-1}^{2}+f x_{l}^{2}$. So there is a non-singular $P$ such that

$$
\begin{aligned}
& P M P^{\mathrm{T}}=\left(\begin{array}{llll}
d_{1} & & & 0 \\
& \ddots & & \\
& & d_{l-1} & \\
0 & & & f
\end{array}\right) \\
& \operatorname{det}(P)^{2} \operatorname{det}(M)=f \prod_{i=1}^{l-1} d_{i}
\end{aligned}
$$

We assume $\operatorname{det}(M)>0$, so $f \prod_{i=1}^{l-1} d_{i}>0$, so $f>0$. Thus, the form is positive definite.

We consider graphs with the following properties:
(i) the graph is connected;
(ii) any two distinct vertices are joined by $0,1,2$ or 3 edges;
(iii) the associated quadratic form is positive definite.

The Dynkin diagram of a semisimple Lie algebra has connected components satisfying (i)-(iii). It is possible to determine all graphs satisfying (i)-(iii).

Theorem 10.6. The only graphs satisfying (i)-(iii) are


Proof. The given graphs clearly satisfy (i) and (ii). We show that they satisfy (iii). We show $Q\left(x_{1}, \ldots, x_{l}\right)$ is positive definite by induction on $l$. If $l=1$ we have $Q\left(x_{1}\right)=2 x_{1}^{2}$, which is positive definite. Suppose $l>1$. There is a vertex $l$ such that when it is removed we have another graph on the list. By induction, $Q\left(x_{1}, \ldots, x_{l-1}\right)$ is positive definite, so all leading minors of the matrix of $Q\left(x_{1}, \ldots, x_{l-1}\right)$ have positive determinant. To complete the induction we show that the matrix of $Q\left(x_{1}, \ldots, x_{l}\right)$ has positive determinant.

Let $Y_{l}$ be a graph of $l$ vertices and $y_{l}$ the determinant of the matrix of the associated quadratic form. In the case $l=1, a_{1}=|2|=2$. In the case $l=2$ we have

$$
\begin{gathered}
a_{2}=\left|\begin{array}{cc}
2 & -1 \\
-1 & 2
\end{array}\right|=3 \\
b_{2}=\left|\begin{array}{cc}
2 & -\sqrt{2} \\
-\sqrt{2} & 2
\end{array}\right|=2 \\
c_{2}=\left|\begin{array}{cc}
2 & -\sqrt{3} \\
-\sqrt{3} & 2
\end{array}\right|=1
\end{gathered}
$$

Suppose $l \geq 3$. Remove a vertex $l$ joined to just one other vertex $l-1$ by a single edge. If $Y_{l}$ is the given graph, let $Y_{l-1}$ be the graph with vertex $l$ removed in this way, and let $Y_{l-2}$ be the graph with vertices $l$ and $l-1$ removed in this way.

$$
y_{l}=\operatorname{det}\left(Y_{l}\right)=\left|\begin{array}{cccccc} 
& & & & & 0 \\
& & & & & \\
& & \\
& & & & 2 & -1 \\
0 & \cdots & 0 & -1 & 2
\end{array}\right|=2 y_{l-1}-(-1)(-1) y_{l-2}=2 y_{l-1}-y_{l-2}
$$

Hence:

Type $A_{l}$

$$
a_{l}=2 a_{l-1}-a_{l-2} \Rightarrow a_{l}=l+1
$$

Type $B_{l}$

$$
b_{l}=2 b_{l-1}-b_{l-2} \Rightarrow b_{l}=2
$$

$$
d_{4}=2 a_{3}-a_{1}^{2}=4
$$

$$
d_{5}=2 d_{4}-a_{3}=4
$$

$$
\Rightarrow d_{l}=4 \text { by induction }
$$

Type $E_{6}$
$e_{6}=2 d_{5}-a_{4}=3$
Type $E_{7}$
$e_{7}=2 e_{6}-d_{5}=2$
Type $E_{8}$
$e_{8}=2 e_{7}-e_{6}=1$
Type $F_{4}$
Type $G_{2}$

$$
f_{4}=2 b_{3}-a_{2}=1
$$

$$
g_{2}=1
$$

Hence, $Q\left(x_{1}, \ldots, x_{l}\right)$ is positive definite in each case.

In order to show the converse, i.e. that the graphs on our list are the only possible ones, we shall first require some additional results.

Proposition 10.7. For each of the following graphs the corresponding quadratic form $Q\left(x_{1}, \ldots, x_{l}\right)$ has determinant zero.



Proof. In most cases we can calculate the determinant as before, but not in types $\widetilde{A}_{l}, \widetilde{C}_{l}$.
Type $\widetilde{A}_{l}$

$$
\left|\begin{array}{ccccc}
2 & -1 & & & -1 \\
-1 & 2 & -1 & 0 & \\
& -1 & \ddots & & \\
& 0 & & \ddots & -1 \\
-1 & & & -1 & 2
\end{array}\right|=0
$$

since the row
sum is $(0, \ldots, 0)$

Type $\widetilde{C}_{l}$

$$
\left|\begin{array}{ccccc} 
& & & & \\
& & & & \\
& & & & 0 \\
& & & 2 & -\sqrt{2} \\
0 & \cdots & 0 & -\sqrt{2} & 2
\end{array}\right|=2 b_{l}-(-\sqrt{2})(-\sqrt{2}) b_{l-1}=0
$$

Type $\widetilde{B}_{l}$

$$
\widetilde{b}_{3}=2 b_{3}-a_{1}^{2}=2.2-2^{2}=0
$$

Type $\widetilde{D}_{l}$
$\widetilde{d}_{4}=2 d_{4}-a_{1}^{3}=2.4-2^{3}=0$
$\widetilde{e}_{6}=2 e_{6}-a_{5}=2.3-6=0$

Type $\widetilde{E}_{7}$

$$
\tilde{e}_{7}=2 e_{7}-d_{6}=2.2-4=0
$$

Type $\widetilde{E}_{8}$
$\widetilde{e}_{8}=2 e_{8}-e_{7}=2.1-2=0$
Type $\widetilde{F}_{4}$
$\widetilde{f}_{4}=2 f_{4}-b_{3}=2.1-2=0$
Type $\widetilde{G}_{2}$
$\widetilde{g}_{2}=2 g_{2}-a_{1}=2.1-2=0$

Lemma 10.8. Let $Y$ be a graph in which any two vertices are joined by at most three edges. Suppose the corresponding quadratic form is positive definite. Suppose $Y^{\prime}$ is a graph obtained from $Y$ by omitting some of the vertices, or by reducing the number of edges, or both. Then the quadratic form for $Y^{\prime}$ is also positive definite.

We call $Y^{\prime}$ a subgraph of $Y$.

## Example.

$$
Y^{\prime}=\quad \bullet \text { is a subgraph of } Y=
$$

Proof. The quadratic form for $Y$ is

$$
Q\left(x_{1}, \ldots, x_{l}\right)=\sum_{i=1}^{l} 2 x_{i}^{2}-\sum_{1 \leq i \neq j \leq l} \sqrt{n_{i j}} x_{i} x_{j}
$$

The quadratic form for $Y^{\prime}$ is

$$
Q^{\prime}\left(x_{1}, \ldots, x_{m}\right)=\sum_{i=1}^{m} 2 x_{i}^{2}-\sum_{1 \leq i \neq j \leq m} \sqrt{n_{i j}^{\prime}} x_{i} x_{j}
$$

with $m \leq l$ and $n_{i j}^{\prime} \leq n_{i j}$. Suppose $Q^{\prime}$ is not positive definite. Then there exists $\left(y_{1}, \ldots, y_{m}\right) \neq 0$ with $Q^{\prime}\left(y_{1}, \ldots, y_{m}\right) \leq 0$. Consider $Q\left(\left|y_{1}\right|, \ldots,\left|y_{m}\right|, 0, \ldots, 0\right)$. This is

$$
\begin{aligned}
\sum_{i=1}^{m} 2 y_{i}^{2}-\sum_{1 \leq i \neq j \leq m} \sqrt{n_{i j}}\left|y_{i}\right|\left|y_{j}\right| & \leq \sum_{i=1}^{m} 2 y_{i}^{2}-\sum_{1 \leq i i j \leq m} \sqrt{n_{i j}^{\prime}}\left|y_{i}\right|\left|y_{j}\right| \\
& \leq \sum_{i=1}^{m} 2 y_{i}^{2}-\sum_{1 \leq i i j \leq m} \sqrt{n_{i j}} y_{i} y_{j} \\
& \leq Q^{\prime}\left(y_{1}, \ldots, y_{m}\right)
\end{aligned}
$$

So $Q\left(\left|y_{1}\right|, \ldots,\left|y_{m}\right|, 0, \ldots, 0\right) \leq Q^{\prime}\left(y_{1}, \ldots, y_{m}\right) \leq 0$. So $Q\left(x_{1}, \ldots, x_{l}\right)$ is not positive definite, a contradiction.

We now return to the proof of 10.6 .
Suppose $Y$ is some graph satisfying conditions. (i)-(iii). By 10.7 and 10.8 we know that no graph of the form $\widetilde{A}, \widetilde{B}, \widetilde{C}, \widetilde{D}, \widetilde{E}, \widetilde{F}$ or $\widetilde{G}$ can be obtained as a subgraph of $Y$.
(a) $Y$ contains no cycles, for otherwise $Y$ would have a subgraph of the form $\tilde{A}_{l}$.
(b) If $Y$ has a triple edge then $Y=G_{2}$, for otherwise $Y$ would have $\widetilde{G}_{2}$ as a subgraph.
(c) Suppose $Y$ has no triple edge. Then $Y$ can have no more than one double edge, for otherwise $Y$ has a subgraph of type $\widetilde{C}_{l}$.
(d) Suppose $Y$ has one double edge. Then $Y$ has no branch point, for otherwise $Y$ has $\widetilde{B}_{l}$ as a subgraph.
(e) If the double edge is not at one end then $Y=F_{4}$, for otherwise $Y$ has a subgraph $\widetilde{F}_{4}$. If the double edge is at one end, $Y=B_{l}$.
(f) Now suppose $Y$ has only single edges. Then $Y$ cannot have a branch point with four or more branches, for otherwise $Y$ has $\widetilde{D}_{4}$ as a subgraph.
(g) $Y$ can have no more than one branch point, for otherwise $Y$ has a subgraph $\widetilde{D}_{l}$, $l \geq 5$.
(h) If $Y$ has no branch points, $Y=A_{l}$. So suppose $Y$ has just one branch point with three branches of lengths $l_{1} \leq l_{2} \leq l_{3}, l_{1}+l_{2}+l_{3}+1=l$. Then $l_{1}=1$, for otherwise $Y$ would have $\widetilde{E}_{6}$ as a subgraph.
(i) If $l_{1}=l_{2}=1, Y=D_{l}$. Also, $l_{2} \leq 2$, for otherwise $Y$ has $\widetilde{E}_{7}$ as a subgraph.
(j) So assume $l_{1}=1, l_{2}=2$. Then $l_{3} \leq 4$, for otherwise $Y$ has $\widetilde{E}_{8}$ as a subgraph.

$$
\begin{aligned}
& l_{3}=2 \Rightarrow Y=E_{6} \\
& l_{3}=3 \Rightarrow Y=E_{7} \\
& l_{3}=4 \Rightarrow Y=E_{8}
\end{aligned}
$$

Corollary 10.9. Every Dynkin diagram of a semisimple Lie algebra has connected components of type $A_{l}, B_{l}, D_{l}, E_{6}, E_{7}, E_{8}, F_{4}$ and $G_{2}$.

The Cartan matrix $A=\left(a_{i j}\right)$ determines the Dynkin diagram since $n_{i j}=a_{i j} a_{j i}$. However, the Dynkin diagram does not always determine the Cartan matrix. Recall that the $a_{i j}$ satisfy

$$
\begin{gathered}
a_{i i}=2 \\
a_{i j} \in\{0,1,2,3\} \text { for } i \neq j
\end{gathered}
$$

If $n_{i j}=1$ then $a_{i j}=a_{j i}=1$. If $n_{i j}=2$ then $\left(a_{i j}, a_{j i}\right)=(-2,-1)$ or $(-1,-2)$. If $n_{i j}=3$ then $\left(a_{i j}, a_{j i}\right)=(-3,-1)$ or $(-1,-3)$. In this last case the Dynkin diagram is $G_{2}$. We have

$$
A=\left(\begin{array}{cc}
2 & -1 \\
-3 & 2
\end{array}\right) \text { or }\left(\begin{array}{cc}
2 & -3 \\
-1 & 2
\end{array}\right)
$$

and one is obtained from the other by re-labeling the vertices.
Suppose $n_{i j}=2$. If $l=2$ the possibilities for the Cartan matrix are

$$
A=\left(\begin{array}{cc}
2 & -1 \\
-2 & 2
\end{array}\right) \text { or }\left(\begin{array}{cc}
2 & -2 \\
-1 & 2
\end{array}\right)
$$

again obtainable from one another by re-labeling. If $l \geq 3$ there are two possible Cartan matrices:

$$
\begin{gathered}
B_{l}=\left(\begin{array}{ccccc}
2 & -1 & & & \\
-1 & 2 & \ddots & & \\
& \ddots & \ddots & \ddots & \\
& & \ddots & 2 & -1 \\
& & -2 & 2
\end{array}\right) \\
B_{3}=\left(\begin{array}{ccc}
2 & -1 & 0 \\
-1 & 2 & -1 \\
0 & -2 & 2
\end{array}\right) \\
\left\lvert\, \begin{array}{cc}
h_{\alpha_{1}}\left|=\left|h_{\alpha_{2}}\right|=\sqrt{2}\right| h_{\alpha_{3}} \mid \\
\bullet & \bullet
\end{array}\right.
\end{gathered}
$$

$$
\begin{gathered}
C_{l}=\left(\begin{array}{ccccc}
2 & -1 & & & \\
-1 & 2 & \ddots & & \\
& \ddots & \ddots & \ddots & \\
& & \ddots & 2 & -2 \\
& & & -1 & 2
\end{array}\right) \\
C_{3}=\left(\begin{array}{ccc}
2 & -1 & 0 \\
-1 & 2 & -2 \\
0 & -1 & 2
\end{array}\right) \\
\left|h_{\alpha_{1}}\right|=\left|h_{\alpha_{2}}\right|=\frac{1}{\sqrt{2}}\left|h_{\alpha_{3}}\right|
\end{gathered}
$$

We place an arrow on the Dynkin diagram when we have a double or triple edge; the arrow points from the longer root to the shorter one. For example, with $G_{2}$ :


Theorem 10.10. The possible Cartan matrices with connected Dynkin diagrams are (up to permutation of the numbering of the vertices):

$$
\begin{aligned}
& A_{l}=\left(\begin{array}{ccccc}
2 & -1 & & & \\
-1 & 2 & \ddots & & \\
& \ddots & \ddots & \ddots & \\
& & \ddots & 2 & -1 \\
& & & -1 & 2
\end{array}\right) \quad B_{l}=\left(\begin{array}{ccccc}
2 & -1 & & \\
-1 & 2 & \ddots & \\
& \ddots & \ddots & \ddots & \\
& & \ddots & 2 & -1 \\
& & & -2 & 2
\end{array}\right) \\
& C_{l}=\left(\begin{array}{ccccc}
2 & -1 & & & \\
-1 & 2 & \ddots & & \\
& \ddots & \ddots & \ddots & \\
& & \ddots & 2 & -2 \\
& & & -1 & 2
\end{array}\right) \\
& D_{l}=\left(\begin{array}{cccccc}
2 & -1 & & & & \\
-1 & 2 & \ddots & & & \\
& \ddots & \ddots & \ddots & & \\
& & \ddots & 2 & -1 & -1 \\
& & & -1 & 2 & 0 \\
& & & -1 & 0 & 2
\end{array}\right) \\
& E_{6}=\left(\begin{array}{cccccc}
2 & -1 & & & & \\
-1 & 2 & -1 & & & \\
& -1 & 2 & -1 & 0 & -1 \\
& & -1 & 2 & -1 & \\
& & 0 & -1 & 2 & \\
& & -1 & & & 2
\end{array}\right) \\
& E_{7}=\left(\begin{array}{ccccccc}
2 & -1 & & & & & \\
-1 & 2 & -1 & & & & \\
& -1 & 2 & -1 & & & \\
& & -1 & 2 & -1 & 0 & -1 \\
& & & -1 & 2 & -1 & \\
& & & 0 & -1 & 2 & \\
& & & -1 & & & 2
\end{array}\right) \\
& E_{8}=\left(\begin{array}{ccccccccc}
2 & -1 & & & & & & \\
-1 & 2 & -1 & & & & & \\
& -1 & 2 & -1 & & & & \\
& & -1 & 2 & -1 & & & \\
& & & -1 & 2 & -1 & 0 & -1 \\
& & & & -1 & 2 & -1 & \\
& & & & 0 & -1 & 2 & \\
& & & & -1 & & & 2
\end{array}\right) \\
& F_{4}=\left(\begin{array}{cccc}
2 & -1 & 0 & 0 \\
-1 & 2 & -1 & 0 \\
0 & -2 & 2 & -1 \\
0 & 0 & -1 & 2
\end{array}\right) \\
& G_{2}=\left(\begin{array}{cc}
2 & -1 \\
-3 & 2
\end{array}\right)
\end{aligned}
$$

## 11. The Indecomposable Root Systems

A root system is called indecomposable if it has a connected Dynkin diagram.
Case $l=1$. We have only one possibility, $A_{1}$ :


$$
\begin{gathered}
\Phi=\left\{ \pm \alpha_{1}\right\} \\
W=\left\langle s_{\alpha_{1}}\right\rangle ;|W|=2
\end{gathered}
$$

Case $l=2$. Here we have three possibilities


Type $A_{2}$.


Type $B_{2}$.

(In the above diagram, the dashed lines indicating the reflection axes are shown slightly offset for clarity.)

Type $G_{2}$.


$$
\Phi=\left\{ \pm \alpha_{1}, \pm \alpha_{2}, \pm\left(\alpha_{1}+\alpha_{2}\right), \pm\left(2 \alpha_{1}+\alpha_{2}\right), \pm\left(3 \alpha_{1}+\alpha_{2}\right), \pm\left(3 \alpha_{1}+2 \alpha_{2}\right)\right\}
$$

(Again, the dashed lines indicating the reflection axes are shown slightly offset.)
Case $l \geq 3$. Type $A_{l}$. It is convenient to describe the root system of type $A_{l}$ in a Euclidean space of dimension $l+1$.

Let $V$ be an $(l+1)$-dimensional Euclidean space. Let $\left\{\varepsilon_{0}, \ldots, \varepsilon_{n}\right\}$ be an orthogonal basis of vectors of the same length, so $\left\langle\varepsilon_{i}, \varepsilon_{j}\right\rangle=K \delta_{i j}$ for some $K>0$.


Define $h_{\alpha_{1}}=\varepsilon_{0}-\varepsilon_{1}, h_{\alpha_{2}}=\varepsilon_{1}-\varepsilon_{2}, \ldots, h_{\alpha_{l}}=\varepsilon_{l-1}-\varepsilon_{l}$. The $h_{\alpha_{i}}$ are linearly independent.

$$
\begin{gathered}
\left\langle h_{\alpha_{i}}, h_{\alpha_{i}}\right\rangle=2 K \\
\left\langle h_{\alpha_{i}}, h_{\alpha_{j}}\right\rangle=0 \text { if } j \neq i-1, i, i+1 \\
\left\langle h_{\alpha_{i}}, h_{\alpha_{i+1}}\right\rangle=-K \\
a_{i, i+1}=2 \frac{\left\langle h_{\alpha_{i}}, h_{\alpha_{i+1}}\right\rangle}{\left\langle h_{\alpha_{i}}, h_{\alpha_{i}}\right\rangle}=\frac{-2 K}{2 K}=-1
\end{gathered}
$$

Thus for suitable $K$ the $h_{\alpha_{i}}$ form a fundamental system of roots of type $A_{l}$. Let $V_{0}$ be the subspace spanned by these vectors; $\operatorname{dim}\left(V_{0}\right)=l$. Consider the map $V \rightarrow V$ given by $\varepsilon_{0} \leftrightarrow \varepsilon_{1}, \varepsilon_{i} \mapsto \varepsilon_{i}$ for $i \geq 2$.

$$
\begin{gathered}
h_{\alpha_{1}} \mapsto-h_{\alpha_{1}} \\
h_{\alpha_{2}} \mapsto h_{\alpha_{1}}+h_{\alpha_{2}} \\
h_{\alpha_{i}} \mapsto h_{\alpha_{i}} \text { for } i \geq 2
\end{gathered}
$$

This is $s_{\alpha_{1}}$. Similarly, the linear map $V \rightarrow V$ such that $\varepsilon_{i-1} \leftrightarrow \varepsilon_{i}$, all others fixed, is $s_{\alpha_{i}}$.
$W$ is generated by $s_{\alpha_{1}}, \ldots, s_{\alpha_{l}}$. The group generated by all transpositions $\left(\varepsilon_{i-1} \varepsilon_{i}\right)$ is isomorphic to $S_{l+1}$. So we have a homomorphism $S_{l+1} \rightarrow V$. This map is surjective; it is also injective, since any permutation of $\left\{\varepsilon_{0}, \ldots, \varepsilon_{n}\right\}$ that fixes each $h_{\alpha_{i}}$ is the identity. Hence, $W \cong S_{l+1}$.

Each $h_{\alpha}$ has the form $h_{\alpha}=w\left(h_{\alpha_{i}}\right)$ for some $w \in W$ and some $i$. Hence

$$
\Phi=\left\{\varepsilon_{i}-\varepsilon_{j} \mid 0 \leq i \neq j \leq l\right\} .
$$

So $|\Phi|=l(l+1)$.

Type $B_{l}$.


Let $V$ be a Euclidean space of dimension $l$ with basis $\left\{\varepsilon_{1}, \ldots, \varepsilon_{l}\right\}$ such that $\left\langle\varepsilon_{i}, \varepsilon_{j}\right\rangle=K \delta_{i j}$. Define

$$
\begin{aligned}
& h_{\alpha_{1}}=\varepsilon_{1}-\varepsilon_{2} \\
& h_{\alpha_{2}}=\varepsilon_{2}-\varepsilon_{3} \\
& \vdots \\
& h_{\alpha_{l-1}}=\varepsilon_{l-1}-\varepsilon_{l} \\
& h_{\alpha_{l}}=\varepsilon_{l}
\end{aligned}
$$

These form a fundamental system of vectors of type $B_{l}$.

$$
\begin{gathered}
\left|h_{\alpha_{1}}\right|=\ldots=\left|h_{\alpha_{l-1}}\right|=\sqrt{2}\left|h_{\alpha_{l}}\right| \\
\left\langle h_{\alpha_{i}}, h_{\alpha_{l}}\right\rangle=0 \text { for } 1 \leq i \leq l-2 \\
\left\langle h_{\alpha_{l-1}}, h_{\alpha_{l}}\right\rangle=-K \\
2 \frac{\left\langle h_{\alpha_{l-1}}, h_{\alpha_{l}}\right\rangle}{\left\langle h_{\alpha_{l-1}}, h_{\alpha_{l-1}}\right\rangle}=\frac{-2 K}{2 K}=-1
\end{gathered}
$$

$s_{\alpha_{1}}: \varepsilon_{1} \leftrightarrow \varepsilon_{2}$ and leaves others fixed,
$s_{\alpha_{2}}: \varepsilon_{2} \leftrightarrow \varepsilon_{3}$ and leaves others fixed,
!
$s_{\alpha_{l-1}}: \varepsilon_{l-1} \leftrightarrow \varepsilon_{l}$ and leaves others fixed,
$s_{\alpha_{l}}: \varepsilon_{l} \mapsto-\varepsilon_{l}$ and leaves others fixed.
$W=\left\langle s_{\alpha_{1}}, \ldots, s_{\alpha_{l}}\right\rangle$; for $w \in W, w\left(\varepsilon_{i}\right)= \pm \varepsilon_{j} .|W|=2^{l} l!$. Each $h_{\alpha}$ has the form $h_{\alpha}=w\left(h_{\alpha_{i}}\right)$ for some $w \in W$ and some $i$. Hence,

$$
\Phi=\left\{ \pm \varepsilon_{i} \pm \varepsilon_{j} \mid 1 \leq i \neq j \leq l\right\} \cup\left\{ \pm \varepsilon_{i} \mid 1 \leq i \leq l\right\}
$$

$$
|\Phi|=2^{2}\binom{l}{2}+2 l=2 l(l-1)+2 l=2 l^{2}
$$

Type $C_{l}$.


Let $V$ be a Euclidean space of dimension $l$ with basis $\left\{\varepsilon_{1}, \ldots, \varepsilon_{l}\right\}$ such that $\left\langle\varepsilon_{i}, \varepsilon_{j}\right\rangle=K \delta_{i j}$. Define

$$
\begin{aligned}
& h_{\alpha_{1}}=\varepsilon_{1}-\varepsilon_{2} \\
& \vdots \\
& h_{\alpha_{l-1}}=\varepsilon_{l-1}-\varepsilon_{l} \\
& h_{\alpha_{l}}=2 \varepsilon_{l}
\end{aligned}
$$

$W$ is the same as for $B_{l}$. For $w \in W, w\left(\varepsilon_{i}\right)= \pm \varepsilon_{j} .|W|=2^{l} l$ !. The $h_{\alpha}$ are the vectors $\pm \varepsilon_{i} \pm \varepsilon_{j}$ (for $i \neq j$ ) and $\pm 2 \varepsilon_{i}$. Hence,

$$
|\Phi|=2^{2}\binom{l}{2}+2 l=2 l(l-1)+2 l=2 l^{2}
$$

Type $D_{l}$.


Let $V$ be a Euclidean space of dimension $l$ with basis $\left\{\varepsilon_{1}, \ldots, \varepsilon_{l}\right\}$ such that $\left\langle\varepsilon_{i}, \varepsilon_{j}\right\rangle=K \delta_{i j}$. Define

$$
\begin{gathered}
h_{\alpha_{1}}=\varepsilon_{1}-\varepsilon_{2} \\
\vdots \\
h_{\alpha_{l-1}}=\varepsilon_{l-1}-\varepsilon_{l} \\
h_{\alpha_{l}}=\varepsilon_{l-1}+\varepsilon_{l}
\end{gathered}
$$

This is a fundamental system of roots of type $D_{l}$.

$$
\begin{gathered}
s_{\alpha_{1}}: \varepsilon_{1} \leftrightarrow \varepsilon_{2} \text { and leaves others fixed, } \\
\vdots \\
s_{\alpha_{l-1}}: \varepsilon_{l-1} \leftrightarrow \varepsilon_{l} \text { and leaves others fixed, } \\
s_{\alpha_{l}}: \pm \varepsilon_{l-1} \mapsto \mp \varepsilon_{l} \text { and leaves others fixed. }
\end{gathered}
$$

For $w \in W, w\left(\varepsilon_{i}\right)= \pm \varepsilon_{j}$. There will be an even number of sign changes, so $|W|=2^{l-1} l$ !. The $h_{\alpha}$ have the form $\pm \varepsilon_{i} \pm \varepsilon_{j}$ for $i \neq j$. So

$$
|\Phi|=2^{2}\binom{l}{2}=4 \frac{l(l-1)}{2}=2 l(l-1)
$$

Type $F_{4}$.


Let $V$ be a 4-dimensional Euclidean space with basis $\left\{\varepsilon_{1}, \varepsilon_{2}, \varepsilon_{3}, \varepsilon_{4}\right\},\left\langle\varepsilon_{i}, \varepsilon_{j}\right\rangle=K \delta_{i j}$.

$$
\begin{gathered}
h_{\alpha_{1}}=\varepsilon_{1}-\varepsilon_{2} \\
h_{\alpha_{2}}=\varepsilon_{2}-\varepsilon_{3} \\
h_{\alpha_{3}}=\varepsilon_{3} \\
h_{\alpha_{4}}=\frac{1}{2}\left(-\varepsilon_{1}-\varepsilon_{2}-\varepsilon_{3}+\varepsilon_{4}\right)
\end{gathered}
$$

This is a fundamental system of vectors of type $F_{4} . s_{\alpha_{1}}, s_{\alpha_{2}}, s_{\alpha_{3}}$ permute $\varepsilon_{1}, \varepsilon_{2}, \varepsilon_{3}$ and change signs arbitrarily.

$$
s_{\alpha_{4}}:\left\{\begin{array}{l}
\varepsilon_{1} \mapsto \frac{1}{2}\left(\varepsilon_{1}-\varepsilon_{2}-\varepsilon_{3}+\varepsilon_{4}\right)=\varepsilon_{1}+h_{\alpha_{4}} \\
\varepsilon_{2} \mapsto \frac{1}{2}\left(-\varepsilon_{1}+\varepsilon_{2}-\varepsilon_{3}+\varepsilon_{4}\right)=\varepsilon_{2}+h_{\alpha_{4}} \\
\varepsilon_{3} \mapsto \frac{1}{2}\left(-\varepsilon_{1}-\varepsilon_{2}+\varepsilon_{3}+\varepsilon_{4}\right)=\varepsilon_{3}+h_{\alpha_{4}} \\
\varepsilon_{4} \leftrightarrow \frac{1}{2}\left(\varepsilon_{1}+\varepsilon_{2}+\varepsilon_{3}+\varepsilon_{4}\right) \\
\varepsilon_{1}+\varepsilon_{2} \mapsto-\varepsilon_{3}+\varepsilon_{4}
\end{array}\right.
$$

Let $S=\left\{h_{\alpha} \mid \alpha \in \Phi\right\}$.

$$
\begin{gathered}
\left\{ \pm \varepsilon_{i} \pm \varepsilon_{j} \mid 1 \leq i \neq j \leq 3\right\} \subseteq S \\
\left\{ \pm \varepsilon_{i} \mid 1 \leq i \leq 3\right\} \subseteq S
\end{gathered}
$$

$$
\begin{gathered}
\frac{1}{2}\left( \pm \varepsilon_{1} \pm \varepsilon_{2} \pm \varepsilon_{3} \pm \varepsilon_{4}\right) \in S \\
\left\{ \pm \varepsilon_{i} \pm \varepsilon_{4} \mid 1 \leq i \leq 3\right\} \subseteq S \\
\pm \varepsilon_{4} \in S
\end{gathered}
$$

So $S$ contains

$$
\begin{gathered}
\pm \varepsilon_{i} \pm \varepsilon_{j} \text { for } 1 \leq i \neq j \leq 4 \\
\pm \varepsilon_{i} \text { for } 1 \leq i \leq 4 \\
\frac{1}{2}\left( \pm \varepsilon_{1} \pm \varepsilon_{2} \pm \varepsilon_{3} \pm \varepsilon_{4}\right)
\end{gathered}
$$

This collection of vectors is closed under the actions of $s_{\alpha_{1}}, \ldots, s_{\alpha_{4}}$.

$$
|\Phi|=2^{2}\binom{4}{2}+2.4+2^{4}=48
$$

(We have 24 short roots and 24 long ones.)
Type $E_{8}$.


Let $V$ be a Euclidean space of dimension 8 with basis $\left\{\varepsilon_{1}, \ldots, \varepsilon_{8}\right\}$ such that $\left\langle\varepsilon_{i}, \varepsilon_{j}\right\rangle=K \delta_{i j}$. Define

$$
\begin{gathered}
\left.\begin{array}{c}
h_{\alpha_{1}}=\varepsilon_{1}-\varepsilon_{2} \\
\vdots \\
h_{\alpha_{6}}=\varepsilon_{6}-\varepsilon_{7} \\
h_{\alpha_{7}}=\varepsilon_{6}+\varepsilon_{7}
\end{array}\right\} \text { type } D_{7} \\
h_{\alpha_{8}}=-\frac{1}{2}\left(\varepsilon_{1}+\varepsilon_{2}+\varepsilon_{3}+\varepsilon_{4}+\varepsilon_{5}+\varepsilon_{6}+\varepsilon_{7}+\varepsilon_{8}\right) \\
\left|h_{a_{7}}\right|=\left|h_{\alpha_{8}}\right| \\
\left\langle h_{\alpha_{7}}, h_{\alpha_{8}}\right\rangle=-K \\
2 \frac{\left\langle h_{\alpha_{7}}, h_{\alpha_{8}}\right\rangle}{\left\langle h_{\alpha_{7}}, h_{\alpha_{7}}\right\rangle}=\frac{-2 K}{2 K}=-1
\end{gathered}
$$

These vectors form a fundamental system of type $E_{8} . s_{\alpha_{1}}, \ldots, s_{\alpha_{7}}$ permute $\varepsilon_{1}, \ldots, \varepsilon_{7}$ and change an even number of signs.

$$
s_{\alpha_{8}}: \varepsilon_{i} \mapsto \frac{1}{4}\left(-\varepsilon_{1}+\ldots+3 \varepsilon_{i}+\ldots-\varepsilon_{8}\right)=\varepsilon_{i}+\frac{1}{2} h_{\alpha_{8}}
$$

Let $S=\left\{h_{\alpha} \mid \alpha \in \Phi\right\} . S$ contains $\pm \varepsilon_{i} \pm \varepsilon_{j}$ for $1 \leq i \neq j \leq 7 . S$ contains $\frac{1}{2} \sum_{i=1}^{8} \varepsilon_{i}$ so $S$ contains $\frac{1}{2}\left(-\varepsilon_{1}-\ldots-\varepsilon_{6}+\varepsilon_{7}+\varepsilon_{8}\right)=s_{\alpha_{8}}\left(\varepsilon_{7}+\varepsilon_{8}\right)$. So $\varepsilon_{7}+\varepsilon_{8} \in S$. So $S$ contains $\pm \varepsilon_{i} \pm \varepsilon_{8}$ for $1 \leq i \leq 7 . S$ also contains $\frac{1}{2}\left( \pm \varepsilon_{1} \pm \ldots \pm \varepsilon_{8}\right)$ with an even number of negative signs. So $S$ contains

$$
\begin{gathered}
\pm \varepsilon_{i} \pm \varepsilon_{j} \text { for } 1 \leq i \neq j \leq 8 \\
\frac{1}{2}\left( \pm \varepsilon_{1} \pm \ldots \pm \varepsilon_{8}\right) \text { with an even number of negative signs }
\end{gathered}
$$

This is the whole of $S$, for it is invariant under $s_{\alpha_{1}}, \ldots, s_{\alpha_{8}}$ : this is clear for $s_{\alpha_{1}}, \ldots, s_{\alpha_{7}}$ but requires a little work to check for $s_{\alpha_{8}}$. So the roots of $E_{8}$ are

$$
\begin{gathered}
\pm \varepsilon_{i} \pm \varepsilon_{j} \text { for } 1 \leq i \neq j \leq 8 \\
\frac{1}{2} \sum_{i=1}^{8} \pm \varepsilon_{i} \text { where } \Pi( \pm)=1 \\
|\Phi|=2^{2}\binom{8}{2}+2^{7}=240
\end{gathered}
$$

Type $E_{7}$.


Take $V$ as before $-\operatorname{dim}(V)=8$. Take $V_{0}$ to be the subspace of $V$ perpendicular to $\varepsilon_{1}-\varepsilon_{8} . \operatorname{dim}\left(V_{0}\right)=7$ and $h_{\alpha_{2}}, \ldots, h_{\alpha_{8}}$ form a basis of $V_{0}$. This is a fundamental system of type $E_{7}$. Consider $S=\left\{h_{\alpha} \mid \alpha \in \Phi\left(E_{7}\right)\right\}$. This set lies in $\left\{h_{\alpha} \mid \alpha \in \Phi\left(E_{8}\right)\right\} \cap V_{0}$. This intersection is

$$
\begin{gathered}
\pm \varepsilon_{i} \pm \varepsilon_{j} \text { for } 2 \leq i \neq j \leq 7 \\
\pm\left(\varepsilon_{1}+\varepsilon_{8}\right) \\
\frac{1}{2}\left(\varepsilon_{1} \pm \varepsilon_{2} \pm \ldots \pm \varepsilon_{7}+\varepsilon_{8}\right) \text { where } \Pi( \pm)=1 \\
\frac{1}{2}\left(-\varepsilon_{1} \pm \varepsilon_{2} \pm \ldots \pm \varepsilon_{7}-\varepsilon_{8}\right) \text { where } \Pi( \pm)=1
\end{gathered}
$$

All of these can be obtained from $h_{\alpha_{2}}, \ldots, h_{\alpha_{8}}$ by means of $s_{\alpha_{2}}, \ldots, s_{\alpha_{8}}$. This is obvious except for $\pm\left(\varepsilon_{1}+\varepsilon_{8}\right)$.

$$
s_{\alpha_{8}}: \varepsilon_{1}+\varepsilon_{2} \leftrightarrow \frac{1}{2}\left(\varepsilon_{1}-\varepsilon_{2}-\ldots-\varepsilon_{7}+\varepsilon_{8}\right)
$$

So $\pm\left(\varepsilon_{1}+\varepsilon_{8}\right) \in S$. So $S$ is

$$
\begin{gathered}
\pm \varepsilon_{i} \pm \varepsilon_{j} \text { for } 2 \leq i \neq j \leq 7 \\
\pm\left(\varepsilon_{1}+\varepsilon_{8}\right) \\
\frac{1}{2}\left(\varepsilon_{1} \pm \varepsilon_{2} \pm \ldots \pm \varepsilon_{7}+\varepsilon_{8}\right) \text { where } \Pi( \pm)=1 \\
\frac{1}{2}\left(-\varepsilon_{1} \pm \varepsilon_{2} \pm \ldots \pm \varepsilon_{7}-\varepsilon_{8}\right) \text { where } \Pi( \pm)=1 \\
\left|\Phi\left(E_{7}\right)\right|=2^{2}\binom{6}{2}+2+2^{5}+2^{5}
\end{gathered}
$$

Type $E_{6}$.


We proceed as before. $h_{\alpha_{3}}, \ldots, h_{\alpha_{8}}$ form a fundamental system of vectors of type $E_{6}$. Let $V_{0}$ be the subspace of $V$ for $E_{8}$ that is orthogonal to $\varepsilon_{1}-\varepsilon_{8}$ and $\varepsilon_{2}-\varepsilon_{8}$. $\operatorname{dim}\left(V_{0}\right)=6$ and $h_{\alpha_{3}}, \ldots, h_{\alpha_{8}}$ form a basis of this space.

$$
\left\{h_{\alpha} \mid \alpha \in \Phi\left(E_{6}\right)\right\} \subseteq\left\{h_{\alpha} \mid \alpha \in \Phi\left(E_{8}\right)\right\} \cap V_{0}
$$

The $h_{\alpha}$ in $V_{0}$ are

$$
\begin{gathered}
\pm \varepsilon_{i} \pm \varepsilon_{j} \text { for } 3 \leq i \neq j \leq 7 \\
\frac{1}{2}\left(\varepsilon_{1}+\varepsilon_{2} \pm \varepsilon_{3} \pm \ldots \pm \varepsilon_{7}+\varepsilon_{8}\right) \text { where } \Pi( \pm)=1 \\
\frac{1}{2}\left(-\varepsilon_{1}-\varepsilon_{2} \pm \varepsilon_{3} \pm \ldots \pm \varepsilon_{7}-\varepsilon_{8}\right) \text { where } \Pi( \pm)=1
\end{gathered}
$$

All of these are obtainable from $h_{\alpha_{3}}, \ldots, h_{\alpha_{8}}$ by $s_{\alpha_{3}}, \ldots, s_{\alpha_{8}}$.

$$
\left|\Phi\left(E_{6}\right)\right|=2^{2}\binom{5}{2}+2^{4}+2^{4}=72
$$

Theorem 11.1. The number of roots in each of the indecomposable root systems is

$$
\begin{array}{ccccccccc}
A_{l} & B_{l} & C_{l} & D_{l} & E_{6} & E_{7} & E_{8} & F_{4} & G_{2} \\
l(l+2) & l(2 l+1) & l(2 l+1) & l(2 l-1) & 78 & 133 & 248 & 52 & 14
\end{array}
$$

## 12. The Semisimple Lie Algebras

Theorem 12.1. (a) If a semisimple Lie algebra $L$ has connected Dynkin diagram $\Delta$ then $L$ is simple.
(b) If $L$ is a semisimple Lie algebra whose Dynkin diagram $\Delta$ has connected components $\Delta_{1}, \ldots, \Delta_{r}$ then $L=L_{1} \oplus \cdots \oplus L_{r}$ where $L_{i}$ is a simple Lie algebra with Dynkin diagram $\Delta_{i}$.

Proof. (a) Let $L=H \oplus\left(\oplus_{\alpha \in \Phi} L_{\alpha}\right)$ be a Cartan decomposition with connected Dynkin diagram $\Delta$. Let $0 \neq I \triangleleft L$. We first show that $I \cap H \neq 0$. Suppose not, i.e. that $I \cap H=0$. Let $0 \neq x \in I$ with $x=h+\sum_{\alpha} \mu_{\alpha} e_{\alpha}$ and the number of non-zero $\mu_{\alpha}$ as small as possible. Let $\mu_{\beta} \neq 0$.

$$
\left[x h_{\beta}\right]=\sum_{\alpha} \mu_{\alpha}\left[e_{\alpha} h_{\beta}\right]=\sum_{\alpha} \mu_{\alpha} \alpha\left(h_{\beta}\right) e_{\alpha}
$$

By 8.7 we can choose $e_{-\beta}$ with $\left[e_{-\beta} e_{\beta}\right]=h_{\beta}$.

$$
\left[\left[x h_{\beta}\right] e_{-\beta}\right]=-\mu_{\beta} \beta\left(h_{\beta}\right) h_{\beta}+\sum_{\substack{\alpha \in \Phi \\ \alpha-\beta \in \Phi}} \mu_{\alpha} \alpha\left(h_{\beta}\right) N_{\alpha,-\beta} e_{\alpha-\beta}
$$

$\left[\left[x h_{\beta}\right] e_{-\beta}\right] \in I$ is non-zero since

$$
-\mu_{\beta} \beta\left(h_{\beta}\right) h_{\beta}=-\underbrace{\mu_{\beta}}_{\neq 0} \underbrace{\left\langle h_{\beta}, h_{\beta}\right\rangle}_{\neq 0} h_{\beta} \neq 0
$$

The number of non-zero $\mu_{\alpha}$ with $\alpha-\beta \in \Phi$ is less than before, a contradiction. Hence, $I \cap H \neq 0$.

We next show that $I \supseteq H$. Suppose not. Then $0 \subset I \cap H \subset H . I \cap H$ is not orthogonal to all $h_{\alpha_{i}}, \alpha_{i} \in \Pi$. For suppose $I \cap H$ is not orthogonal to $h_{\alpha_{i}}$. Let $x \in I \cap H$ be such that $\left\langle x, h_{\alpha_{i}}\right\rangle \neq 0$. Then

$$
\left[e_{\alpha_{i}} x\right]=\alpha_{i}(x) e_{\alpha_{i}}=\left\langle h_{\alpha_{i}}, x\right\rangle e_{\alpha_{i}} \in I
$$

So $e_{\alpha_{i}} \in I$. So $\left[e_{-\alpha_{i}} e_{\alpha_{i}}\right]=h_{\alpha_{i}} \in I$. So for each $\alpha_{i} \in \Pi$ either $\left\langle I \cap H, h_{\alpha_{i}}\right\rangle=0$ or $h_{\alpha_{i}} \in I$. Both classes are non-empty. Choose $h_{\alpha_{j}} \notin I$; then $\left\langle I \cap H, h_{\alpha_{j}}\right\rangle=0$. This means $\Delta$ is disconnected, a contradiction. Hence, $H \subseteq I$.

Now let $\alpha \in \Phi$.

$$
\left[e_{\alpha} h_{\alpha}\right]=\alpha\left(h_{\alpha}\right) e_{\alpha}=\underbrace{\left\langle h_{\alpha}, h_{\alpha}\right\rangle}_{\neq 0} e_{\alpha}
$$

So $e_{\alpha} \in I$. Hence $I$ contains $H$ and all $e_{\alpha}$. So $I=L ; L$ is simple.
(b) Suppose $\Delta$ is the disjoint union of connected components $\Delta_{1}, \ldots, \Delta_{r}$. Then $\Pi$ is the union of orthogonal components $\Pi_{1}, \ldots, \Pi_{r}$. Let $H_{i}$ be the subspace spanned by $\left\{h_{\alpha} \mid \alpha \in \Pi_{i}\right\}$. Then $H=H_{1} \oplus \ldots \oplus H_{r}$ and the $H_{i}$ are mutually orthogonal. Now consider $s_{\alpha}$ for some $\alpha \in \Pi_{i}$. Then $\alpha$ transforms $H_{i}$ into itself and fixes each vector in $H_{j}$ for $j \neq i ; s_{\alpha}\left(H_{j}\right)=H_{j}$. Since the $s_{\alpha}$ for $\alpha \in \Pi$ generate $W, w\left(H_{j}\right)=H_{j}$ for each $w \in W$.

For each $\alpha \in \Phi, h_{\alpha}=w\left(h_{\alpha_{i}}\right)$ for some $w \in W$ and some $i$. So $h_{\alpha} \in H_{i}$ for some $i$. Let $\Phi_{i}=\left\{\alpha \in \Phi \mid h_{\alpha} \in H_{i}\right\}$. Then $\Phi=\Phi_{1} \cup \ldots \cup \Phi_{r}$. Let $L_{i}$ be the subspace of $L$ spanned by $H_{i}$ and the $L_{\alpha}$ with $\alpha \in \Phi_{i}$. We see that $L=L_{1} \oplus \ldots \oplus L_{r}$ as a direct sum of vector spaces.

To see that $L_{i}$ is a subalgebra of $L$ it is sufficient to show that $\alpha, \beta \in \Phi_{i} \Rightarrow\left[e_{\alpha} e_{\beta}\right] \in L_{i}$. If $\alpha+\beta \notin \Phi$ then $\left[e_{\alpha} e_{\beta}\right]=0$. If $\beta=-\alpha$ then $\left[e_{\alpha} e_{-\alpha}\right]=-h_{\alpha} \in H_{i}$. If $\alpha+\beta \in \Phi$ then $\alpha+\beta \in \Phi_{i}$ and $h_{\alpha+\beta}=h_{\alpha}+h_{\beta} \in H_{i}$. So $L_{i}$ is a subalgebra.

We next check that $i \neq j \Rightarrow\left[L_{i} L_{j}\right]=0$. Let $\alpha \in \Phi_{i}, \beta \in \Phi_{j}$.

$$
\begin{gathered}
{\left[h_{\alpha} h_{\beta}\right]=0} \\
{\left[h_{\alpha} e_{\beta}\right]=0 \text { since }\left\langle h_{\alpha}, h_{\beta}\right\rangle=0} \\
{\left[e_{\alpha} h_{\beta}\right]=0 \text { since }\left\langle h_{\alpha}, h_{\beta}\right\rangle=0} \\
{\left[e_{\alpha} e_{\beta}\right]=0 \text { since } \alpha+\beta \notin \Phi}
\end{gathered}
$$

$h_{\alpha}+h_{\beta}$ does not lie in any $H_{k}, h_{\alpha} \in H_{i}, h_{\beta} \in H_{j}$. So $\left[L_{i} L_{j}\right]=0$ for $i \neq j$ :

$$
\left[x_{1}+\ldots+x_{r}, y_{1}+\ldots+y_{r}\right]=\left[x_{1} y_{1}\right]+\ldots+\left[x_{r} y_{r}\right]
$$

So $L=L_{1} \oplus \ldots \oplus L_{r}$ as a direct sum of Lie algebras. We now see that each $L_{i}$ is semisimple. Let $I \triangleleft L_{i}$ be soluble. $\left[I L_{j}\right]=0$ if $i \neq j$, so $I \triangleleft L . I$ is a soluble ideal of $L$, but $L$ is semisimple, so $I=0$, so $L_{i}$ is semisimple.

We now show that $H_{i}$ is a Cartan subalgebra of $L_{i} . H$ is a Cartan subalgebra of $L$, so there is a regular element $x \in L$ such that $H$ is the 0 -(generalized) eigenspace of ad $x$. Since $x \in H$ and $H=H_{1} \oplus \ldots \oplus H_{r}$ we can write $x=x_{1}+\ldots+x_{r}$ with $x_{i} \in H_{i}$.

$$
\operatorname{ad} x_{i}: L_{j} \rightarrow\left\{\begin{array}{rl}
L_{j} & i=j \\
0 & i \neq j
\end{array}\right.
$$

So the 0 -eigenspace of $\operatorname{ad} x$ on $L$ is the direct sum of the 0 -eigenspaces of the $\mathrm{ad} x_{i}$ on the $L_{i} . x$ is regular in $L$ if and only if each $x_{i}$ is regular in $L_{i}$. So each $x_{i}$ is regular in $L_{i}$ and the 0-eigenspace of $\operatorname{ad} x_{i}$ in $L_{i}$ is $H_{i}$. So $H_{i}$ is a Cartan subalgebra of $L_{i}$.

$$
L_{i}=H_{i} \oplus\left(\oplus_{\alpha \in \Phi_{i}} L_{\alpha}\right)
$$

is a Cartan decomposition of $L_{i}$. So $\Phi_{i}$ is the root system of $L_{i} ; \Pi_{i}$ is a fundamental root system of $L_{i}$; the Dynkin diagram of $L_{i}$ is $\Delta_{i}$. But $\Delta_{i}$ is connected, so $L_{i}$ is simple by (a).

We next consider simple Lie algebras with a given indecomposable Cartan matrix $A$.
Existence Problem: Is there a simple Lie algebra with given Cartan matrix $A$ ?
Isomorphism Problem: Are any two such Lie algebras isomorphic?
Let $L$ be a simple Lie algebra and $H$ a Cartan subalgebra of $L$ :

$$
\begin{gathered}
L=H \oplus\left(\oplus_{\alpha \in \Phi} L_{\alpha}\right) \\
\Phi=\Phi^{+} \cup \Phi^{-}
\end{gathered}
$$

For each $\alpha \in \Phi^{+}$choose $0 \neq e_{\alpha} \in L_{\alpha} ; \quad L_{\alpha}=\mathbb{C} e_{\alpha}$. Choose $e_{-\alpha} \in L_{-\alpha}$ such that $\left[e_{-\alpha} e_{\alpha}\right]=h_{\alpha}$. If $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$ then the $h_{\alpha_{i}}$ and $e_{\alpha}$ form a basis of $L .\left[L_{\alpha} L_{\beta}\right] \subseteq L_{\alpha+\beta}$ so $\left[e_{\alpha} e_{\beta}\right]=N_{\alpha, \beta} e_{\alpha+\beta}$ if $\alpha+\beta \in \Phi, \alpha \neq-\beta$.

$$
\begin{gathered}
{\left[h_{\alpha_{i}} h_{\alpha_{j}}\right]=0} \\
{\left[e_{\alpha} h_{\alpha_{i}}\right]=\left\langle h_{\alpha}, h_{\alpha_{i}}\right) e_{\alpha}} \\
{\left[e_{-\alpha} e_{\alpha}\right]=h_{\alpha}} \\
{\left[e_{\alpha} e_{\beta}\right]=\left\{\begin{array}{cc}
N_{\alpha, \beta} e_{\alpha+\beta} & \alpha+\beta \in \Phi \\
0 & \text { otherwise }
\end{array}\right.}
\end{gathered}
$$

The $N_{\alpha, \beta}$ are called the structure constants.

Proposition 12.2. The structure constants $N_{\alpha, \beta}$ satisfy
(i) $N_{\alpha, \beta}=-N_{\beta, \alpha}$;
(ii) if $\alpha, \beta, \gamma \in \Phi$ have $\alpha+\beta+\gamma=0$ then $N_{\alpha, \beta}=N_{\beta, \gamma}=N_{\gamma, \alpha}$;
(ii) if $\alpha, \beta, \gamma, \delta \in \Phi$ have $\alpha+\beta+\gamma+\delta=0$ and no pair have sum zero then

$$
N_{\alpha, \beta} N_{\gamma, \delta}+N_{\beta, \gamma} N_{\alpha, \delta}+N_{\gamma, \alpha} N_{\beta, \delta}=0
$$

(if $\xi, \eta \in \Phi, \eta \neq-\xi, \eta+\xi \notin \Phi$ take $N_{\xi, \eta}=0$ );
(iv) if $\alpha, \beta \in \Phi$ have $\alpha+\beta \in \Phi$ then

$$
N_{\alpha, \beta} N_{-\alpha,-\beta}=-\frac{(p+1) q}{2}\left\langle h_{\alpha}, h_{\alpha}\right\rangle
$$

where the $\alpha$-chain of roots through $\beta$ is

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

In particular, $N_{\alpha, \beta} \neq 0$, so $\left[L_{\alpha} L_{\beta}\right]=L_{\alpha+\beta}$.

Proof. (i) $\left[e_{\alpha} e_{\beta}\right]=-\left[e_{\beta} e_{\alpha}\right]$ so $N_{\alpha, \beta}=-N_{\beta, \alpha}$.
(ii) Suppose $\alpha+\beta+\gamma=0$.

$$
\begin{array}{cc} 
& {\left[\left[e_{\alpha} e_{\beta}\right] e_{\gamma}\right]+\left[\left[e_{\beta} e_{\gamma}\right] e_{\alpha}\right]+\left[\left[e_{\gamma} e_{\alpha}\right] e_{\beta}\right]=0} \\
\Rightarrow & N_{\alpha, \beta}\left[e_{\alpha+\beta} e_{\gamma}\right]+N_{\beta, \gamma}\left[e_{\beta+\gamma} e_{\alpha}\right]+N_{\gamma, \alpha}\left[e_{\gamma+\alpha} e_{\beta}\right]=0 \\
\Rightarrow & N_{\alpha, \beta} h_{\gamma}+N_{\beta, \gamma} h_{\alpha}+N_{\gamma, \alpha} h_{\beta}=0 \\
\Rightarrow & \left(-N_{\alpha, \beta}+N_{\beta, \gamma}\right) h_{\alpha}+\left(-N_{\alpha, \beta}+N_{\gamma, \alpha}\right) h_{\beta}=0
\end{array}
$$

Since $h_{\alpha}$ and $h_{\beta}$ are linearly independent, $N_{\alpha, \beta}=N_{\beta, \gamma}=N_{\gamma, \alpha}$.
(iii) Take $\alpha, \beta, \gamma, \delta \in \Phi$ with zero sum and no opposite pairs.

$$
\begin{array}{cc} 
& {\left[\left[e_{\alpha} e_{\beta}\right] e_{\gamma}\right]+\left[\left[e_{\beta} e_{\gamma}\right] e_{\alpha}\right]+\left[\left[e_{\gamma} e_{\alpha}\right] e_{\beta}\right]=0} \\
\Rightarrow & N_{\alpha, \beta}\left[e_{\alpha+\beta} e_{\gamma}\right]+N_{\beta, \gamma}\left[e_{\beta+\gamma} e_{\alpha}\right]+N_{\gamma, \alpha}\left[e_{\gamma+\alpha} e_{\beta}\right]=0 \\
\Rightarrow & \left(N_{\alpha, \beta} N_{\alpha+\beta, \gamma}+N_{\beta, \gamma} N_{\beta+\gamma, \alpha}+N_{\gamma, \alpha} N_{\gamma+\alpha, \beta}\right) e_{\alpha+\beta+\gamma}=0 \\
\Rightarrow & N_{\alpha, \beta} N_{\gamma, \delta}+N_{\beta, \gamma} N_{\alpha, \delta}+N_{\gamma, \alpha} N_{\beta, \delta}=0
\end{array}
$$

(iv) Let $\alpha, \beta \in \Phi$ with $\alpha+\beta \in \Phi$.

$$
\begin{array}{lc} 
& {\left[\left[e_{\alpha} e_{-\alpha}\right] e_{\beta}\right]+\left[\left[e_{-\alpha} e_{\beta}\right] e_{\alpha}\right]+\left[\left[e_{\beta} e_{\alpha}\right] e_{-\alpha}\right]=0} \\
\Rightarrow & -\left[h_{\alpha} e_{\beta}\right]+N_{-\alpha, \beta}\left[e_{-\alpha+\beta} e_{\alpha}\right]+N_{\beta, \alpha}\left[e_{\alpha+\beta} e_{-\alpha}\right]=0 \\
\Rightarrow & \left(\beta\left(h_{\alpha}\right)+N_{-\alpha, \beta} N_{-\alpha+\beta, \alpha}+N_{\beta, \alpha} N_{\alpha+\beta,-\alpha}\right) e_{\beta}=0 \\
\Rightarrow & N_{\alpha, \beta} N_{-\alpha+\beta, \alpha}+N_{\beta, \alpha} N_{\alpha+\beta,-\alpha}=-\left\langle h_{\alpha}, h_{\beta}\right\rangle \\
\Rightarrow & N_{\alpha, \beta} N_{-\alpha,-\beta}-N_{\alpha,-\alpha+\beta} N_{-\alpha, \alpha-\beta}=\left\langle h_{\alpha}, h_{\beta}\right\rangle
\end{array}
$$

Take $M_{\alpha, \beta}=N_{\alpha, \beta} N_{-\alpha,-\beta}$, so

$$
M_{\alpha, \beta}-M_{\alpha,-\alpha+\beta}=\left\langle h_{\alpha}, h_{\beta}\right\rangle
$$

Let the $\alpha$-chain of roots through $\beta$ be

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

So

$$
\begin{aligned}
M_{\alpha, \beta}-M_{\alpha,-\alpha+\beta} & =\left\langle h_{\alpha}, h_{\beta}\right\rangle \\
M_{\alpha,-\alpha+\beta}-M_{\alpha,-2 \alpha+\beta} & =\left\langle h_{\alpha}, h_{-\alpha+\beta}\right\rangle \\
M_{\alpha,-2 \alpha+\beta}-M_{\alpha,-3 \alpha+\beta} & =\left\langle h_{\alpha}, h_{-2 \alpha+\beta}\right\rangle \\
& \vdots \\
M_{\alpha,-p \alpha+\beta} & =\left\langle h_{\alpha}, h_{-p \alpha+\beta}\right\rangle
\end{aligned}
$$

So

$$
\begin{aligned}
M_{\alpha, \beta} & =(p+1)\left\langle h_{\alpha}, h_{\beta}\right\rangle-\left\langle h_{\alpha}, h_{\alpha}\right\rangle(1+2+\ldots+p) \\
& =(p+1)\left\langle h_{\alpha}, h_{\beta}\right\rangle-\frac{p(p+1)}{2}\left\langle h_{\alpha}, h_{\alpha}\right\rangle
\end{aligned}
$$

By, by $8.11,2\left\langle h_{\alpha}, h_{\beta}\right\rangle /\left\langle h_{\alpha}, h_{\alpha}\right\rangle=p-q$, so

$$
\begin{aligned}
M_{\alpha, \beta} & =\left\langle h_{\alpha}, h_{\alpha}\right\rangle\left(\frac{(p+1)(p-q)}{2}-\frac{p(p+1)}{2}\right) \\
& =-\frac{(p+1) q}{2}\left\langle h_{\alpha}, h_{\alpha}\right\rangle
\end{aligned}
$$

So $N_{\alpha, \beta} \neq 0 ;\left[L_{\alpha} L_{\beta}\right]=L_{\alpha+\beta}$.

This result has certain consequences. Let $\alpha, \beta \in \Phi, \alpha+\beta \in \Phi,\left[e_{\alpha} e_{\beta}\right]=N_{\alpha, \beta} e_{\alpha+\beta}$. Let $\gamma=-\alpha-\beta$. Then $\alpha+\beta+\gamma=0$. We have the following ordered pairs of roots whose sum is a root:

$$
\begin{array}{cccccc}
(\alpha, \beta) & (\beta, \gamma) & (\gamma, \alpha) & (\beta, \alpha) & (\gamma, \beta) & (\alpha, \gamma) \\
(-\alpha,-\beta) & (-\beta,-\gamma) & (-\gamma,-\alpha) & (-\beta,-\alpha) & (-\gamma,-\beta) & (-\alpha,-\gamma)
\end{array}
$$

We have a total order $\alpha \prec \beta$. An ordered pair $(\alpha, \beta)$ such that $0 \prec \alpha \prec \beta$ is called a special pair.

Either one or two of $\alpha, \beta, \gamma$ are positive; if one is positive two of $-\alpha,-\beta,-\gamma$ are positive. Of the twelve pairs above just one is special.
$N_{\alpha, \beta}$, for any ordered pair $(\alpha, \beta)$, can be expressed in terms of $N_{\xi, \eta}$ for $(\xi, \eta)$ a special pair by using 12.2(i), (ii), (iv). So consider $N_{\alpha, \beta}$ when $(\alpha, \beta)$ is special; $\alpha+\beta \in \Phi^{+} \backslash \Pi$. This root may be expressible as $\alpha+\beta=\alpha^{\prime}+\beta^{\prime}$ where $\left(\alpha^{\prime}, \beta^{\prime}\right)$ is special and distinct from $(\alpha, \beta)$.

A special pair $(\alpha, \beta)$ is called extra special if for any special pair $\left(\alpha^{\prime}, \beta^{\prime}\right)$ with $\alpha+\beta=\alpha^{\prime}+\beta^{\prime}$ we have $\alpha \preceq \alpha^{\prime}$.

The number of extra special pairs is $\left|\Phi^{+} \backslash \Pi\right|$.
Now let $\left(\alpha^{\prime}, \beta^{\prime}\right)$ be special but not extra special. Then $\alpha+\beta=\alpha^{\prime}+\beta^{\prime}$ where $(\alpha, \beta)$ is extra special - such an extra special pair exists because the set of special and extra special pairs is finite.

$$
\begin{gathered}
\alpha^{\prime}+\beta^{\prime}+(-\alpha)+(-\beta)=0 \\
N_{\alpha^{\prime}, \beta^{\prime}} N_{-\alpha,-\beta}+N_{\beta^{\prime},-\alpha} N_{\alpha^{\prime},-\beta}+N_{-\alpha, \alpha^{\prime}} N_{\beta^{\prime},-\beta}=0 \\
0 \prec \alpha \prec \alpha^{\prime} \prec \beta^{\prime} \prec \beta \\
N_{\alpha^{\prime}, \beta^{\prime}} N_{-\alpha,-\beta}+N_{\beta-\alpha^{\prime}, \alpha^{\prime}} N_{-\alpha,-\left(\beta^{\prime}-\alpha\right)}+N_{\beta-\beta^{\prime}, \beta^{\prime}} N_{-\left(\alpha^{\prime}-\alpha\right),-\alpha}=0
\end{gathered}
$$

We show that $N_{\alpha^{\prime}, \beta^{\prime}}$ is determined by $N_{\xi, \eta}$ 's for extra special pairs $(\xi, \eta)$. We use induction on $\alpha^{\prime}+\beta^{\prime}$ :
$N_{\alpha^{\prime}, \beta^{\prime}}$ is determined by $N_{\alpha, \beta}, N_{\beta-\alpha^{\prime}, \alpha^{\prime}}, N_{\alpha, \beta^{\prime}-\alpha}, N_{\beta-\beta^{\prime}, \beta^{\prime}}, N_{\alpha^{\prime}-\alpha, \alpha} .(\alpha, \beta)$ is extra special.

Either $\left(\beta-\alpha^{\prime}, \alpha^{\prime}\right)$ or $\left(\alpha^{\prime}, \beta-\alpha^{\prime}\right)$ is special:

$$
\left(\beta-\alpha^{\prime}\right)+\alpha^{\prime}=\beta \prec \alpha+\beta=\alpha^{\prime}+\beta^{\prime}
$$

So $N_{\beta-\alpha^{\prime}, \alpha^{\prime}}$ can be expressed in terms of $N_{\xi, \eta}$ for extra special pairs $(\xi, \eta)$.
Either $\left(\alpha, \beta^{\prime}-\alpha\right)$ or $\left(\beta^{\prime}-\alpha, \alpha\right)$ is special:

$$
\alpha+\left(\beta^{\prime}-\alpha\right)=\beta^{\prime} \prec \alpha^{\prime}+\beta^{\prime}
$$

Either $\left(\beta-\beta^{\prime}, \beta^{\prime}\right)$ or $\left(\beta^{\prime}, \beta-\beta^{\prime}\right)$ is special:

$$
\left(\beta-\beta^{\prime}\right)+\beta^{\prime}=\beta \prec \alpha^{\prime}+\beta^{\prime}
$$

Either $\left(\alpha^{\prime}-\alpha, \alpha\right)$ or $\left(\alpha, \alpha^{\prime}-\alpha\right)$ is special:

$$
\left(\alpha^{\prime}-\alpha\right)+\alpha=\alpha^{\prime} \prec \alpha^{\prime}+\beta^{\prime}
$$

Hence, $N_{\alpha^{\prime}, \beta^{\prime}}$ can be expressed in terms of $N_{\xi, \eta}$ 's for extra special pairs $(\xi, \eta)$. So relations 12.2(i)-(iv) expresses all $N_{\alpha^{\prime}, \beta^{\prime}}$ 's in terms of $N_{\xi, \eta}$ 's for extra special pairs $(\xi, \eta)$.

Theorem 12.3. There is a unique simple Lie algebra, up to isomorphism, with a given indecomposable Cartan matrix.

$$
\begin{gathered}
L=H \oplus\left(\oplus_{\alpha \in \Phi} L_{\alpha}\right) \\
\operatorname{dim}(L)=l+|\Phi|
\end{gathered}
$$

Thus, the simple Lie algebras and their dimensions are given by

$$
\begin{array}{cccc}
\operatorname{dim}\left(A_{l}\right)=l(l+2) & \operatorname{dim}\left(B_{l}\right)=l(2 l+1) & \operatorname{dim}\left(C_{l}\right)=l(2 l+1) & \operatorname{dim}\left(D_{l}\right)=l(2 l-1) \\
(l \geq 1) & (l \geq 2) & (l \geq 3) & (l \geq 4) \\
\operatorname{dim}\left(E_{6}\right)=78 & \operatorname{dim}\left(E_{7}\right)=133 & \operatorname{dim}\left(E_{8}\right)=248 \quad \operatorname{dim}\left(F_{4}\right)=52 & \operatorname{dim}\left(G_{2}\right)=14
\end{array}
$$

Note. The following are isomorphic:

$$
\begin{gathered}
B_{2} \cong C_{2} \\
A_{3} \cong D_{3} \\
D_{2} \cong A_{1} \oplus A_{1}
\end{gathered}
$$



Proof. Uniqueness. Let $L, L^{\prime}$ be simple Lie algebras with indecomposable Cartan matrix $A=\left(a_{i j}\right) . L$ has a Cartan decomposition $L=H \oplus\left(\oplus_{\alpha \in \Phi} L_{\alpha}\right)$. If $\Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$ then $H$
has basis $\left\{h_{\alpha_{1}}, \ldots, h_{\alpha_{l}}\right\}$; $L$ has basis $\left\{h_{\alpha_{1}}, \ldots, h_{\alpha_{l}}\right\} \cup\left\{e_{\alpha} \mid \alpha \in \Phi\right\}$. Multiplication of basis elements:

$$
\begin{gathered}
{\left[h_{\alpha_{i}} h_{\alpha_{j}}\right]=0} \\
{\left[e_{\alpha} h_{\alpha_{i}}\right]=\left\langle h_{\alpha}, h_{\alpha_{i}}\right) e_{\alpha}} \\
{\left[e_{-\alpha} e_{\alpha}\right]=h_{\alpha}} \\
{\left[e_{\alpha} e_{\beta}\right]=\left\{\begin{array}{cc}
N_{\alpha, \beta} e_{\alpha+\beta} & \alpha+\beta \in \Phi \\
0 & 0 \neq \alpha+\beta \notin \Phi
\end{array}\right.}
\end{gathered}
$$

All scalar products $\left\langle h_{\alpha}, h_{\beta}\right\rangle$ are determined by $A$. Also, all of the $h_{\alpha}$ (as linear combinations of the $h_{\alpha_{i}}$ ) are determined by $A$.

$$
s_{\alpha_{i}}\left(h_{\alpha_{j}}\right)=h_{\alpha_{j}}-a_{i j} h_{\alpha_{i}}
$$

So the $s_{\alpha_{i}}$ are determined by $A$. The Weyl group $W$ is generated by $s_{\alpha_{1}}, \ldots, s_{\alpha_{l}}$. So $W$ is determined by $A . h_{\alpha}=w\left(h_{\alpha_{i}}\right)$ for some $i$ and some $w \in W$. So the $h_{\alpha}$ are determined by $A$.

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

So $2\left\langle h_{\alpha}, h_{\beta}\right\rangle /\left\langle h_{\alpha}, h_{\alpha}\right\rangle$ is determined by $A$. But

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

by 8.13. So $\left\langle h_{\alpha}, h_{\alpha}\right\rangle$ is determined by $A$. So $\left\langle h_{\alpha}, h_{\beta}\right\rangle$ is determined by $A$.
Suppose a basis $\left\{h_{\alpha_{1}}^{\prime}, \ldots, h_{\alpha_{l}}^{\prime}\right\} \cup\left\{e_{\alpha}^{\prime} \mid \alpha \in \Phi\right\}$ of $L^{\prime}$ is given. We describe how to choose a basis of $L$. The $h_{\alpha_{i}}$ are uniquely determined. Choose $e_{\alpha} \neq 0$ in $L_{\alpha}$ for each $\alpha \in \Pi$. For each $\alpha \in \Phi^{+} \backslash \Pi$ there is a unique extra special pair $(\beta, \gamma)$ such that $\alpha=\beta+\gamma$, $\beta, \gamma \prec \alpha$.

Assume by induction that $e_{\beta}, e_{\gamma}$ are already chosen. Choose $e_{\alpha}$ by $e_{\alpha}=N_{\beta, \gamma}\left[e_{\beta} e_{\gamma}\right]$ where $N_{\beta, \gamma}=N_{\beta, \gamma}^{\prime}$, the structure constant for $L^{\prime}$. Having thus chosen $e_{\alpha}$ for $\alpha \in \Phi^{+}$, we choose $e_{-\alpha}$ by $\left[e_{-\alpha} e_{\alpha}\right]=h_{\alpha}$.

The $N_{\alpha, \beta}$ for arbitrary $\alpha, \beta$ are determined by the $N_{\xi, \eta}$, where $(\xi, \eta)$ is extra special, by 12.2. Since $N_{\xi, \eta}=N_{\xi, \eta}^{\prime}$ for all extra special $(\xi, \eta)$ it follows that $N_{\alpha, \beta}=N_{\alpha, \beta}^{\prime}$ for all $\alpha, \beta \in \Phi$ with $\alpha+\beta \in \Phi$.

This shows that $L$ and $L^{\prime}$ are isomorphic.
Existence. (Sketch proof.) Begin with Cartan matrix $A=\left(a_{i j}\right)$. Let $H$ be an $l$ dimensional vector space over $\mathbb{C}$ with basis $h_{\alpha_{1}}, \ldots, h_{\alpha_{l}}$. We define $s_{\alpha_{i}}: H \rightarrow H$ by $s_{\alpha_{i}}\left(h_{\alpha_{j}}\right)=h_{\alpha_{j}}-a_{i j} h_{\alpha_{i}}$, a self-inverse map. Let $W$ be the group of all non-singular linear maps $H \rightarrow H$ generated by $s_{\alpha_{1}}, \ldots, s_{\alpha_{l}}$. $W$ is finite. Correspondingly,

$$
\left\{h_{\alpha}=w\left(h_{\alpha_{i}}\right) \mid w \in W, 1 \leq i \leq l\right\}
$$

is also finite. (The $h_{\alpha}$ were determined in Chapter 11.) We now define a bilinear map

$$
\begin{gathered}
H \times H \rightarrow \mathbb{C} \\
(x, y) \mapsto\langle x, y\rangle
\end{gathered}
$$

This form is uniquely determined by $A$. Define $\alpha \in H^{*}$ by $\alpha(x)=\left\langle h_{\alpha}, x\right\rangle$; let $\Phi$ be the set of all such $\alpha$.

Let $L$ be a vector space over $\mathbb{C}$ with $\operatorname{dim}(L)=\operatorname{dim}(H)+|\Phi|$ with basis

$$
\left\{h_{\alpha_{1}}, \ldots, h_{\alpha_{l}}\right\} \cup\left\{e_{\alpha} \mid \alpha \in \Phi\right\}
$$

Define a bilinear map

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

We define [ ] on the basis elements by

$$
\begin{gathered}
{\left[h_{\alpha_{i}} h_{\alpha_{j}}\right]=0} \\
{\left[e_{\alpha} h_{\alpha_{i}}\right]=-\left[h_{\alpha_{i}} e_{\alpha}\right]=\left\langle h_{\alpha}, h_{\alpha_{i}}\right\rangle e_{\alpha}} \\
{\left[e_{\alpha} e_{\beta}\right]=\left\{\begin{array}{cc}
\left.e_{-\alpha} e_{\alpha}\right]=h_{\alpha} \\
e_{\alpha+\beta} & \alpha+\beta \in \Phi \\
0 & 0 \neq \alpha+\beta \notin \Phi
\end{array}\right.}
\end{gathered}
$$

The $N_{\alpha, \beta}$ can be chosen arbitrarily if $(\alpha, \beta)$ is extra special, e.g. $N_{\alpha, \beta}=1 . N_{\alpha, \beta}$ is determined for all other pairs by 12.2. So multiplication of basis elements is determined by $A$. We make various checks:

Check $[x x]=0$ for all $x \in L$. (Easy.) Check $[[x y] z]+[[y z] x]+[[z x] y]+0$. (Most are easy, but $x=e_{\alpha}, y=e_{\beta}, z=e_{\gamma}$ is difficult.) Then $L$ is a Lie algebra, $L=H \oplus\left(\oplus_{\alpha \in \Phi} L_{\alpha}\right)$, $L_{\alpha}=\mathbb{C} e_{\alpha}$. Check that $H$ is a Cartan subalgebra of $L$. (Difficult.) Then $L=H \oplus\left(\oplus_{\alpha \in \Phi} L_{\alpha}\right)$ is a Cartan decomposition of $L$ with respect to $H$. (Easy.) Then $\Phi$ is the set of roots of $L$ with respect to $H . \Pi=\left\{\alpha_{1}, \ldots, \alpha_{l}\right\}$ is a fundamental system of roots inside $\Phi$. We have

$$
2 \frac{\left\langle h_{\alpha_{i}}, h_{\alpha_{j}}\right\rangle}{\left\langle h_{\alpha_{i}}, h_{\alpha_{i}}\right\rangle}=a_{i j}
$$

so $A$ is the Cartan matrix. Finally, the argument of 12.1(a) proves that $L$ is simple.

## Review.



If we choose a different Cartan subalgebra and a different fundamental system do we get a different $A$ ?

Theorem 12.4. (i) Let L be a Lie algebra and $H_{1}, H_{2}$ Cartan subalgebras. Then there exists an automorphism $\theta: L \rightarrow L$ such that $\theta\left(H_{1}\right)=H_{2}$.
(ii) A subalgebra $H$ of $L$ is a Cartan subalgebra if and only if $H$ is nilpotent and $H=\mathcal{N}(H)$.

Theorem 12.5. Let $\Phi$ be the root system of a semisimple Lie algebra and let $\Pi_{1}, \Pi_{2}$ be two fundamental systems in $\Phi$. Then there is $a$ w $\in W$ such that $w\left(\Pi_{1}\right)=\Pi_{2}$.
12.4 and 12.5 imply that the Cartan matrix is uniquely determined by $L$. So the simple Lie algebras on our list are pairwise non-isomorphic.

We have four infinite families of simple Lie algebras and five exceptional ones:

## Classical

## Exceptional

| $A_{l}$ | $B_{l}$ | $C_{l}$ | $D_{l}$ | $E_{6}$ | $E_{7}$ | $E_{8}$ | $F_{4}$ | $G_{2}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $l(l+2)$ | $l(2 l+1)$ | $l(2 l+1)$ | $l(2 l-1)$ | 78 | 133 | 248 | 52 | 14 |
| $l \geq 1$ | $l \geq 2$ | $l \geq 3$ | $l \geq 4$ |  |  |  |  |  |

Type $A_{l}$. We can write $\operatorname{dim}\left(A_{l}\right)=l(l+2)=(l+1)^{2}-1$. The set $\mathfrak{s l}_{l+1}(\mathbb{C})$ of all $(l+1) \times(l+1)$ matrices of trace zero forms a Lie algebra of type $A_{l}$. The diagonal subalgebra is a Cartan subalgebra.

Type $B_{l}$. The set $\mathfrak{s o}_{2 l+1}(\mathbb{C})$ of all $(2 l+1) \times(2 l+1)$ matrices $X$ satisfying

$$
X^{\mathrm{T}}\left(\begin{array}{lll}
1 & 0 & 0 \\
0 & 0 & I_{l} \\
0 & I_{l} & 0
\end{array}\right)=-\left(\begin{array}{ccc}
1 & 0 & 0 \\
0 & 0 & I_{l} \\
0 & I_{l} & 0
\end{array}\right) X
$$

forms a simple Lie algebra of type $B_{l}$. The diagonal subalgebra is a Cartan subalgebra. $\mathfrak{s o}_{2 l+1}(\mathbb{C})$ is isomorphic to the Lie algebra of all $(2 l+1) \times(2 l+1)$ skew-symmetric matrices. Elements of $\mathfrak{s o}_{2 l+1}(\mathbb{C})$ have the block form

$$
X=\left(\begin{array}{ccc}
0 & X_{01} & X_{02} \\
-X_{02}^{\mathrm{T}} & X_{11} & X_{12} \\
-X_{01}^{\mathrm{T}} & X_{22} & -X_{11}^{\mathrm{T}}
\end{array}\right)
$$

where $X_{11}$ is an arbitrary $l \times l$ matrix, $X_{12}$ and $X_{21}$ are $l \times l$ symmetric matrices and $X_{01}$ and $X_{02}$ are arbitrary $1 \times l$ matrices (row vectors).

Type $C_{l}$. The set $\mathfrak{s p}_{2 l}(\mathbb{C})$ of all $2 l \times 2 l$ matrices $X$ satisfying

$$
X^{\mathrm{T}}\left(\begin{array}{cc}
0 & I_{l} \\
-I_{l} & 0
\end{array}\right)=-\left(\begin{array}{cc}
0 & I_{l} \\
-I_{l} & 0
\end{array}\right) X
$$

forms a simple Lie algebra of type $C_{l}$. The diagonal subalgebra is a Cartan subalgebra. Elements of $\mathfrak{s p}_{2 l}(\mathbb{C})$ have the block form

$$
X=\left(\begin{array}{cc}
X_{11} & X_{12} \\
X_{21} & -X_{11}^{\mathrm{T}}
\end{array}\right)
$$

where $X_{11}$ is an arbitrary $l \times l$ matrix and $X_{12}$ and $X_{21}$ are $l \times l$ symmetric matrices.

Type $D_{l}$. The set $\mathfrak{s o}_{2 l}(\mathbb{C})$ of all $2 l \times 2 l$ matrices $X$ such that

$$
X^{\mathrm{T}}\left(\begin{array}{cc}
0 & I_{l} \\
I_{l} & 0
\end{array}\right)=-\left(\begin{array}{cc}
0 & I_{l} \\
I_{l} & 0
\end{array}\right) X
$$

forms a simple Lie algebra of type $D_{l}$. The diagonal subalgebra is a Cartan subalgebra. $\mathfrak{s o}_{2 l}(\mathbb{C})$ is isomorphic to the Lie algebra of all $2 l \times 2 l$ skew-symmetric matrices. Elements of $\mathfrak{s o}_{2 l}(\mathbb{C})$ have the block form

$$
X=\left(\begin{array}{cc}
X_{11} & X_{12} \\
X_{21} & -X_{11}^{\mathrm{T}}
\end{array}\right)
$$

where $X_{11}$ is an arbitrary $l \times l$ matrix and $X_{12}$ and $X_{21}$ are $l \times l$ skew-symmetric matrices.
$\mathfrak{s l}_{m}(\mathbb{C})$ is the Lie algebra of $\operatorname{SL}_{m}(\mathbb{C})=\left\{X \in \mathrm{GL}_{m}(\mathbb{C}) \mid \operatorname{det}(X)=1\right\} ; \mathfrak{s o}_{m}(\mathbb{C})$ is the Lie algebra of $\mathrm{SO}_{m}(\mathbb{C})=\left\{X \in \mathrm{GL}_{m}(\mathbb{C}) \mid X^{\mathrm{T}} X=I_{m}\right.$ and $\left.\operatorname{det}(X)=1\right\}$.

Type $G_{2} . \operatorname{dim}\left(G_{2}\right)=14$. Consider the algebra of octonians (a.k.a. Cayley numbers), $\mathcal{O}$. $\operatorname{dim}(\mathcal{O})=8 . \mathcal{O}$ has basis $1, e_{1}, e_{2}, \ldots, e_{7}:$

1 is the multiplicative identity;

$$
\begin{gathered}
e_{i}^{2}=-1 \text { for } 1 \leq i \leq 7 ; \\
e_{i} e_{j}= \pm e_{k} \text { for } 1 \leq i \neq j \leq 7 .
\end{gathered}
$$



The projective plane over the 2 -field.

$$
e_{i} e_{j}=e_{k} \text { if } i \rightarrow j ; e_{i} e_{j}=-e_{k} \text { if } i \leftarrow j
$$

$\mathcal{O}$ is a non-associative algebra. The set of all derivations of $\mathcal{O}$, i.e. linear maps $D: \mathcal{O} \rightarrow \mathcal{O}$ such that $D(x y)=D(x) y+x D(y)$, forms a Lie algebra of type $G_{2}$.

Type $F_{4}$. Define the octonian conjugate:

$$
\begin{gathered}
x=a_{0} 1+\sum_{i=1}^{7} a_{i} e_{i} \\
\bar{x}=a_{0} 1-\sum_{i=1}^{7} a_{i} e_{i} \\
x=\bar{x} \Leftrightarrow x=a_{0} 1
\end{gathered}
$$

A matrix $M$ over $\mathcal{O}$ is called Hermitian if $M^{\mathrm{T}}=\bar{M}$. Let $\mathcal{J}$ be the $\mathbb{C}$-vector space of all $3 \times 3$ Hermitian matrices over $\mathcal{O}$. Such matrices have the form

$$
M=\left(\begin{array}{ccc}
a 1 & x & y \\
\bar{x} & b 1 & z \\
\bar{y} & \bar{z} & c 1
\end{array}\right)
$$

where $a, b, c \in \mathbb{C}$ and $x, y, z \in \mathcal{O} \cdot \operatorname{dim}(\mathcal{J})=27$. We define multiplication on $\mathcal{J}$ by

$$
\begin{aligned}
& M_{1} \times M_{2}=\frac{1}{2}\left(M_{1} M_{2}+M_{2} M_{1}\right) \\
& M_{1} \times M_{2} \in \mathcal{J} \text { for } M_{1}, M_{2} \in \mathcal{J}
\end{aligned}
$$

$\mathcal{J}$ is a commutative non-associative algebra; it is an example of a Jordan algebra, the axioms for which are that

$$
\begin{aligned}
X \times Y & =Y \times X \\
\left(X^{2} \times Y\right) \times X & =X^{2} \times(Y \times X)
\end{aligned}
$$

The derivations of $\mathcal{J}$ form a simple Lie algebra of type $F_{4}$.
$E_{6}, E_{7}$ and $E_{8}$ can all be described in terms of $\mathcal{O}$ and $\mathcal{J}$.
There is an alternative approach to the existence theorem, which proceeds (in outline) as follows:

Let $L$ be a simple Lie algebra with Cartan matrix $A=\left(a_{i j}\right), L=H \oplus\left(\oplus_{\alpha \in \Phi} L_{\alpha}\right) . H$ has basis $h_{\alpha_{1}}, \ldots, h_{\alpha_{l}}$. Let

$$
h_{i}=\frac{2 h_{\alpha_{i}}}{\left\langle h_{\alpha_{i}}, h_{\alpha_{i}}\right\rangle} .
$$

$h_{1}, \ldots, h_{l}$ also form a basis of $H$. Choose $0 \neq e_{i} \in L_{\alpha_{i}}$ and $0 \neq f_{i} \in L_{-\alpha_{i}}$.

$$
\begin{aligned}
& {\left[e_{j} h_{i}\right] }=\alpha_{j}\left(h_{i}\right) e_{j} \\
&=\left\langle h_{\alpha_{j}}, h_{i}\right\rangle e_{j} \\
&\left.=\frac{2\left\langle h_{\alpha_{i}}, h_{\alpha_{j}}\right\rangle}{\left\langle h_{\alpha_{i}}, h_{\alpha_{i}}\right\rangle}\right\rangle \\
&=a_{i j} e_{j} \\
& {\left[f_{j} h_{i}\right]=-a_{i j} f_{j} }
\end{aligned}
$$

Choose $f_{i}$ with $\left[f_{i} e_{i}\right]=h_{i} ; e_{1}, \ldots, e_{l}$ generate $\oplus_{\alpha \in \mathbb{\Phi}^{+}} L_{\alpha} ; f_{1}, \ldots, f_{l}$ generate $\oplus_{\alpha \in \mathbb{Q}^{-}} L_{\alpha}$; $h_{1}, \ldots, h_{l}$ generate $H$. So $G=\left\{e_{i}, f_{i}, h_{i} \mid 1 \leq i \leq l\right\}$ generates $L$. We have relations $\mathcal{R}$ :

$$
\begin{gathered}
{\left[h_{i} h_{j}\right]=0} \\
{\left[e_{j} h_{i}\right]=a_{i j} e_{j}} \\
{\left[f_{j} h_{i}\right]=-a_{i j} f_{j}} \\
{\left[f_{i} e_{i}\right]=h_{i}} \\
{\left[f_{j} e_{i}\right]=0 \text { if } i \neq j} \\
{\left[e_{i} \ldots e_{i}\left[e_{i} e_{j}\right]\right]=0 \text { if } i \neq j\left(1-a_{i j} e_{i}^{\prime} \text { 's }\right)} \\
{\left[f_{i} \ldots f_{i}\left[f_{i} f_{j}\right]\right]=0 \text { if } i \neq j\left(1-a_{i j} f_{i}^{\prime} \text { s }\right)}
\end{gathered}
$$

(The requirements for $1-a_{i j} e_{i}$ 's and $f_{i}$ 's arise from consideration of the $\alpha_{i}$-chain of roots through $\alpha_{j}$.) The Lie algebra generated by $G$ with relations $\mathcal{R}$ is a finitedimensional Lie algebra with Cartan matrix $A$.
$L$ is constructed as follows: let $R$ be the polynomial ring $\mathbb{C}\left\langle e_{1}, \ldots, e_{l}, f_{1}, \ldots, f_{l}, h_{1}, \ldots, h_{l}\right\rangle$ with non-commutative variables. $[R]$ is the Lie algebra obtained from $R$. Let $M$ be the subalgebra generated by $e_{1}, \ldots, e_{l}, f_{1}, \ldots, f_{l}, h_{1}, \ldots, h_{l}$. Let $I$ be the ideal of $M$ generated by

$$
\left[h_{i} h_{j}\right],\left[e_{j} h_{i}\right]-a_{i j} e_{j},\left[f_{j} h_{i}\right]+a_{i j} f_{j},\left[f_{j} e_{i}\right]-\delta_{i j} h_{i},\left[e_{i} \ldots e_{i}\left[e_{i} e_{j}\right]\right],\left[f_{i} \ldots f_{i}\left[f_{i} f_{j}\right]\right]
$$

Then $L=M / I$. We can show that $L$ is finite-dimensional and has Cartan matrix $A$.


[^0]:    ${ }^{\dagger}$ After Marius Sophus Lie (1842-99).

[^1]:    ${ }^{\dagger}$ The classification of the simple Lie algebras was completed in the 1890's by Élie Cartan and Wilhelm Killing, working independently.

[^2]:    ${ }^{\dagger}$ After Hermann Weyl.

[^3]:    ${ }^{\dagger}$ After E.B. Dynkin. Also due to H.S.M. Coxeter.

