

**TOPICS IN NON COMMUTATIVE ALGEBRAIC  
GEOMETRY AND CONGRUENT ZETA FUNCTIONS  
(PART VI). LIE ALGEBRAS AND THEIR  
ENVELOPING ALGEBRAS**

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1. LIE ALGEBRAS

**DEFINITION 1.1.** Let  $K$  be a commutative ring. Then a **Lie algebra**  $\mathfrak{g}$  over  $K$  is a  $K$ -module with a bilinear (non associative) bracket product (“Lie bracket”)

$$[\bullet, \bullet] : \mathfrak{g} \times \mathfrak{g} \rightarrow \mathfrak{g}$$

which satisfies the following axioms:

- (1)  $[X, X] = 0$  for all  $X \in \mathfrak{g}$ .
- (2) (“Jacobi identity”)

$$[X, [Y, Z]] = [[X, Y], Z] + [Y, [X, Z]] \quad (\forall X, Y, Z \in \mathfrak{g}).$$

**EXAMPLE 1.2.** Any associative algebra  $A$  over  $k$  may be regarded as a Lie algebra with the “commutator” as a Lie bracket.

In this talk, we always regard associative algebra as a Lie algebra equipped with the commutator product unless otherwise specified.

**LEMMA 1.3.** *Let  $\mathfrak{g}$  be a Lie algebra over a ring  $k$ . Then there exists an associative unital algebra  $U(\mathfrak{g})$  with a Lie algebra homomorphism*

$$\iota_{\mathfrak{g}} : \mathfrak{g} \rightarrow U(\mathfrak{g})$$

*with the following universal property:*

*For any associative unital algebra  $A$  with a Lie algebra homomorphism  $\phi : \mathfrak{g} \rightarrow A$ , there exists a unique algebra homomorphism*

$$\psi : U(\mathfrak{g}) \rightarrow A$$

*such that  $\psi \circ \iota_{\mathfrak{g}} = \phi$  holds.*

*The pair  $(U(\mathfrak{g}), \iota_{\mathfrak{g}})$  is unique up to an isomorphism.*

**DEFINITION 1.4.** Under the assumption of the previous Lemma, The pair  $(U(\mathfrak{g}), \iota_{\mathfrak{g}})$  is called the **universal enveloping algebra** the Lie algebra  $\mathfrak{g}$ .

Universal enveloping algebras of Lie algebras form an important class of non commutative associative algebras. Our task in this Part is to describe these algebras in our language.

## 2. REPRESENTATIONS OF A LIE ALGEBRA

DEFINITION 2.1. Let  $k$  be a field. A finite dimensional representation of a Lie algebra  $\mathfrak{g}$  over  $k$  is a Lie algebra homomorphism

$$\rho : \mathfrak{g} \rightarrow M_n(k).$$

Note: The full matrix algebra  $M_n(k)$ , when regarded as a Lie algebra equipped with the commutator product, is commonly denoted as  $\mathfrak{gl}_n(k)$ .

EXAMPLE 2.2. Let  $k$  be a field. Let  $\mathfrak{g}$  be a finite dimensional Lie algebra over  $k$ . We then have an **adjoint representation**

$$\mathfrak{g} \ni X \mapsto \text{ad}(X) = (Y \mapsto [X, Y]) \in \text{End}_{k\text{-linear}}(\mathfrak{g}).$$

## 3. POINCARÉ-BIRKOFF-WITT THEOREM

In this section we prove the Poincaré-Birkoff-Witt Theorem. The treatment here essentially follows [1]. Let  $\mathfrak{g}$  be a Lie algebra over a field  $k$ . To prove the theorem we consider  $S_k(\mathfrak{g})$ , the symmetric algebra of  $\mathfrak{g}$  over  $k$ . Let us denote the multiplication of  $S_k(\mathfrak{g})$  by  $(x, y) \mapsto x \circ y$ . We note that each element  $x$  of  $S_k(\mathfrak{g})$  has its degree  $\text{deg}(x)$ . (as a polynomial in elements of  $\mathfrak{g}$ .)

LEMMA 3.1. *We choose a ordered basis  $(x_\lambda; \lambda \in \Omega)$ . (That means, a basis with a totally ordered index set  $\Omega$ .) Then there exists a linear action of  $\mathfrak{g}$  on  $S_k(\mathfrak{g})$  which obeys the following rules:*

(1) *For any  $x \in \mathfrak{g}$  and for any  $y \in S_k(\mathfrak{g})$ ,*

$$\text{deg}(x.y - x \circ y) \leq \text{deg}(y)$$

(2) *If  $\lambda_0 \leq \lambda_1 \leq \lambda_2 \leq \dots \leq \lambda_n$ , then we have*

$$x_{\lambda_0} \cdot (x_{\lambda_1} \circ x_{\lambda_2} \circ x_{\lambda_3} \circ \dots \circ x_{\lambda_n}) = x_{\lambda_0} \circ x_{\lambda_1} \circ x_{\lambda_2} \circ x_{\lambda_3} \circ \dots \circ x_{\lambda_n}.$$

(3) *For any  $x, y \in \mathfrak{g}$  and for any  $z \in S_k(\mathfrak{g})$ , we have*

$$x.(y.z) - y.(x.z) = [x, y].z$$

The proof is done by a careful use of induction. Namely,

SUBLEMMA 3.2. *We employ the same assumption of the above Lemma. Then for each  $m \in \mathbb{Z}_{>0}$ , there exists a unique  $k$ -bilinear map*

$$f_m : \mathfrak{g} \times S_k(\mathfrak{g})_{\leq m} \rightarrow S_k(\mathfrak{g})_{\leq m+1}$$

*which obeys the following rules:*

(1) For any  $x \in \mathfrak{g}$  and for any  $y \in S_k(\mathfrak{g})_{\leq m}$ ,

$$\deg(f_m(x, y) - x \circ y) \leq \deg(y)$$

(2) If  $\lambda_0 \leq \lambda_1 \leq \lambda_2 \leq \cdots \leq \lambda_n$  and  $n \leq m$ , then we have

$$f_m(x_{\lambda_0}, x_{\lambda_1} \circ x_{\lambda_2} \circ x_{\lambda_3} \circ \cdots \circ x_{\lambda_n}) = x_{\lambda_0} \circ x_{\lambda_1} \circ x_{\lambda_2} \circ x_{\lambda_3} \circ \cdots \circ x_{\lambda_n}.$$

(3) For any  $x, y \in \mathfrak{g}$  and for any  $z \in S_k(\mathfrak{g})_{\leq m-1}$ , we have

$$f_m(x, f_m(y, z)) = f_m(y, f_m(x, z)) + f_m([x, y], z)$$

PROOF. We note first that  $S_k(\mathfrak{g})$  has the set of monomials

$$\{x_{\lambda_1} \circ x_{\lambda_2} \circ x_{\lambda_3} \circ \cdots \circ x_{\lambda_n}; \lambda_1 \leq \lambda_2 \leq \lambda_3 \leq \cdots \leq \lambda_n\}$$

as a  $k$ -basis. For monomial  $w = x_{\lambda_1} \circ x_{\lambda_2} \circ x_{\lambda_3} \circ \cdots \circ x_{\lambda_n}$  such that  $\lambda_1 \leq \lambda_2 \leq \lambda_3 \leq \cdots \leq \lambda_n$ , we put  $z = x_{\lambda_2} \circ x_{\lambda_3} \circ \cdots \circ x_{\lambda_n}$ . Then

$$w = f_{m-1}(x_{\lambda_1}, z).$$

We define inductively the action of  $x_{\lambda_0}$  on it by the following equations.

$$f_m(x_{\lambda_0}, x_{\lambda_1} \circ z) = \begin{cases} x_{\lambda_0} \circ x_{\lambda_1} \circ z & (\text{if } \lambda_0 \leq \lambda_1) \\ \left( \begin{array}{l} x_{\lambda_1} \circ x_{\lambda_0} \circ z \\ + f_{m-1}(x_{\lambda_1}, f_{m-1}(x_{\lambda_0}, z) - x_{\lambda_0} \circ z) \\ + f_{m-1}([x_{\lambda_0}, x_{\lambda_1}], z) \end{array} \right) & (\text{if } \lambda_0 > \lambda_1) \end{cases}$$

We first note that the above definition is necessary to meet our conditions. Indeed, by (2) we necessarily define as above for  $\lambda_0 \leq \lambda_1$ . When  $\lambda_0 > \lambda_1$ , we compute

$$\begin{aligned} & x_{\lambda_0} \cdot (x_{\lambda_1} \circ z) \\ & \stackrel{(3)}{=} x_{\lambda_1} \cdot x_{\lambda_0} \cdot z + [x_{\lambda_0}, x_{\lambda_1}] \cdot z \\ & = x_{\lambda_1} \cdot (x_{\lambda_0} \cdot z - x_{\lambda_0} \circ z) + x_{\lambda_1} \cdot (x_{\lambda_0} \circ z) + [x_{\lambda_0}, x_{\lambda_1}] \cdot z \\ & \stackrel{(2)}{=} x_{\lambda_1} \cdot (x_{\lambda_0} \cdot z - x_{\lambda_0} \circ z) + x_{\lambda_1} \circ x_{\lambda_0} \circ z + [x_{\lambda_0}, x_{\lambda_1}] \cdot z \end{aligned}$$

and take a careful look at degrees of each monomials using (1). From this argument we see in particular that the action is uniquely determined by conditions (1),(2),(3).

It is easy to see that the conditions (1),(2) are satisfied by  $f_m$  defined as above.. Let us proceed to verify that the  $f_m$  so defined also satisfies (3). Let us consider  $x_\lambda, x_\mu, z = x_{\mu_1} \circ x_{\mu_2} \circ \cdots \circ x_{\mu_n}$  with  $\mu_1 \leq \mu_2 \leq \cdots \leq \mu_n$ ,  $n \leq m - 1$ . We need to prove

$$(b) \quad x_\lambda \cdot x_\mu \cdot z - x_\mu \cdot x_\lambda \cdot z = [x_\lambda, x_\mu] \cdot z.$$

Since the equation above is antisymmetric in  $\mu, \nu$ , we may assume that  $\lambda \leq \mu$ .

(i) Case where  $\lambda \leq \mu_1$ .

$$\begin{aligned} & x_\lambda \cdot x_\mu \cdot z \\ &= x_\lambda \cdot (x_\mu \circ z) + x_\lambda \cdot (x_\mu \cdot z - x_\mu \circ z) \\ &\stackrel{(1)}{=} x_\lambda \circ x_\mu \circ z + x_\lambda \cdot (x_\mu \cdot z - x_\mu \circ z) \end{aligned}$$

In other words,

$$f_m(x_\lambda, f_m(x_\mu, z)) = x_\lambda \circ x_\mu \circ z + f_{m-1}(x_\lambda, (f_{m-1}(x_\lambda, z) - x_\mu \circ z)).$$

On the other hand we have

$$\begin{aligned} & x_\mu \cdot x_\lambda \cdot z \\ &= x_\mu \cdot (x_\lambda \circ z) \\ &\stackrel{\text{by def}}{=} x_\lambda \circ x_\mu \circ z + f_{m-1}(x_\lambda, f_{m-1}(x_\mu, z) - x_\mu \circ z) + f_{m-1}([x_\mu, x_\lambda], z) \end{aligned}$$

So the equation  $\flat$  surely holds in this case.

(ii) Case where  $\lambda, \mu > \mu_1$ .

In this case we need to “decompose”  $z$  further:

$$z = x_\nu \cdot w.$$

We first forget about the hypothesis  $\lambda \leq \mu$  and prove

$$\begin{aligned} & x_\lambda \cdot (x_\mu \cdot (x_\nu \cdot w)) \quad (\heartsuit) \\ &= x_\nu \cdot (x_\lambda \cdot (x_\mu \cdot w)) + [x_\lambda, x_\nu] \cdot (x_\mu \cdot w) + [x_\mu, x_\nu] \cdot (x_\lambda \cdot w) + [x_\lambda, [x_\mu, x_\nu]] \cdot w \end{aligned}$$

(Since we are doing induction, we need to pay a special attention on degrees on operands. That means, we should use  $f_m$ 's rather than the above “lazy” notation. But that is fairly cumbersome, so we keep on being lazy here.)

Let us now admit that the above equation  $\heartsuit$  is true and prove the rest of the equation (3). By interchanging  $\lambda$  and  $\mu$  in the equation  $(\heartsuit)$ , we obtain

$$\begin{aligned} & x_\mu \cdot (x_\lambda \cdot (x_\nu \cdot w)) \quad (\diamond) \\ &= x_\nu \cdot (x_\mu \cdot (x_\lambda \cdot w)) + [x_\mu, x_\nu] \cdot (x_\lambda \cdot w) + [x_\lambda, x_\nu] \cdot (x_\mu \cdot w) + [x_\mu, [x_\lambda, x_\nu]] \cdot w \end{aligned}$$

Then by subtracting  $(\diamond)$  from  $(\heartsuit)$ , we obtain

$$\begin{aligned} & x_\lambda \cdot (x_\mu \cdot (x_\nu \cdot w)) - x_\mu \cdot (x_\lambda \cdot (x_\nu \cdot w)) \\ &= x_\nu \cdot (x_\lambda \cdot (x_\mu \cdot w) - x_\mu \cdot (x_\lambda \cdot w)) \\ &\quad + ([x_\lambda, [x_\mu, x_\nu]] - [x_\mu, [x_\lambda, x_\nu]]) \cdot w. \end{aligned}$$

Since  $\deg(w)$  is smaller than  $\deg(z)$ , by induction hypothesis the first term in the right hand side may be replaced by  $x_\nu \cdot ([x_\lambda, x_\mu] \cdot w)$ . The second term may be replaced, by the Jacobian identity, by  $[[x_\lambda, x_\mu], x_\nu] \cdot w$ . So the equation (b) holds in this case too.

It remains to prove the equation ( $\heartsuit$ ). By the induction hypothesis we have

$$x_\mu \cdot (x_\nu \cdot w) = x_\nu \cdot (x_\mu \cdot w) + [x_\mu, x_\nu] \cdot w.$$

Also by the induction hypothesis we have

$$x_\lambda \cdot ([x_\mu, x_\nu] \cdot w) = [x_\mu, x_\nu] \cdot (x_\lambda \cdot w) + [x_\lambda, [x_\mu, x_\nu]] \cdot w$$

Lastly, we decompose  $x_\mu \cdot w$  as

$$x_\mu \cdot w = (x_\mu \circ w) + (x_\mu \cdot w - x_\mu \circ w) = (x_\mu \circ w) + y$$

Then the second term  $y$  has degree smaller than  $\deg(z) = \deg(w) + 1$ . The case (i) applies to the first term and we obtain:

$$x_\lambda \cdot (x_\nu \cdot (x_\mu \cdot w)) = x_\nu \cdot (x_\lambda \cdot (x_\mu \cdot w)) + [x_\lambda, x_\nu] \cdot (x_\mu \cdot w).$$

These altogether complete the proof.  $\square$

**THEOREM 3.3** (Poincaré, Birkoff, Witt(PBW)). *Let  $\mathfrak{g}$  be a Lie algebra over a field  $k$ . Then we have a  $k$ -algebra isomorphism*

$$\Psi : \text{Gr}(U(\mathfrak{g})) \cong S(\mathfrak{g}).$$

**PROOF.** Let

$$\iota_0 : \mathfrak{g} \rightarrow \text{Gr}(U(\mathfrak{g}))$$

be the obvious  $k$ -linear map.

Using the universality of symmetric algebra, there exists a unique  $k$ -algebra homomorphism

$$\Phi : S(\mathfrak{g}) \rightarrow \text{Gr}(U(\mathfrak{g}))$$

which extends  $\iota_0$ . On the other hand the action defined in the Lemma 1.3 gives us a linear map

$$\Psi_0 : U(\mathfrak{g}) \ni x \mapsto x \cdot 1 \in S(\mathfrak{g})$$

which is clearly degree-decreasing. So it defines a  $k$ -linear map

$$\Psi : \text{Gr}(U(\mathfrak{g})) \rightarrow \text{Gr}(S(\mathfrak{g})) \cong S(\mathfrak{g}).$$

Now the composition we obtain

$$\Psi \circ \Phi : S(\mathfrak{g}) \xrightarrow{\Phi} \text{Gr}(U(\mathfrak{g})) \xrightarrow{\Psi} S(\mathfrak{g})$$

coincides with the identity map. Indeed, it coincides with the identity on monomials of the form

$$x_{\lambda_1} \circ x_{\lambda_2} \circ x_{\lambda_3} \circ \cdots \circ x_{\lambda_{n-1}} \circ x_{\lambda_n}.$$

The map  $\Phi$  is easily verified to be surjective. So we conclude that  $\Phi$  and  $\Psi$  are both bijective and are inverse to each other.  $\square$

#### 4. JORDAN-CHEVALLEY DECOMPOSITION OF A SQUARE MATRIX

##### 4.1. Existence and uniqueness of Jordan-Chevalley decomposition.

DEFINITION 4.1. Let  $A$  be a square matrix over a field  $k$ . A **Jordan-Chevalley decomposition** (also called SN-decomposition) of  $A$  is a decomposition of  $A$

$$A = S + N$$

which satisfies the following conditions.

- (1)  $S$  is semisimple (that means, the minimal polynomial of  $S$  has only simple roots.)
- (2)  $N$  is nilpotent.
- (3)  $SN = NS$
- (4)  $S, N \in \overline{k}[A]$

A main objective of this section is to prove the following proposition.

PROPOSITION 4.2. *For any square matrix  $A$  over a field  $k$ , there exists a unique Jordan-Chevalley decomposition.*

To prove it, we need some basic facts from linear algebra.

LEMMA 4.3. *Let  $A$  be a square matrix over a field  $k$ . Let  $m_A(X)$  be the minimal polynomial of  $A$  over  $k$ . If  $m_A$  is decomposed into two coprime polynomials, that means,*

$$m_A = m_1 m_2 \quad (m_1, m_2) = 1,$$

then  $A$  is similar to a direct sum

$$A_1 \oplus A_2 = \begin{pmatrix} A_1 & 0 \\ 0 & A_2 \end{pmatrix}$$

where  $m_1(A_1) = 0$ ,  $m_2(A_2) = 0$ .

PROOF. Since  $k[X]$  is an Euclidean domain, it is a principal ideal domain. thus we see that there exists a polynomial  $l_1(X), l_2(X) \in k[X]$  such that

$$l_1(X)m_1(X) + l_2(X)m_2(X) = 1$$

holds. We put

$$E_j = l_j(A)m_j(A) \quad (j = 1, 2).$$

Then we see easily that  $E_1, E_2$  are mutually orthogonal projection. That means, we have

$$E_1^2 = E_1, E_2^2 = E_2, E_1 + E_2 = 1, E_1 E_2 = 0.$$

It is also easy to see that both  $E_1$  and  $E_2$  commute with  $A$ . Now putting  $A_1 = A|_{\text{Range } E_2}$  and  $A_2 = A|_{\text{Range } E_1}$  we see that

$$A = AE_2 + AE_1 = A_1 \oplus A_2.$$

with  $A_1$  and  $A_2$  satisfying the required property. □

**COROLLARY 4.4.** *Every square matrix  $A$  over a field  $k$  is similar to a direct sum of square matrices  $A_1, A_2, \dots, A_s$  with each minimal polynomial  $m_{A_j}$  equals to a power  $f_j^{e_j}$  of an irreducible polynomial  $f_j$  over  $k$ .*

□

**COROLLARY 4.5.** *When  $k$  is algebraically closed, every square matrix  $A$  over a field  $k$  is similar to a direct sum of square matrices  $A_1, A_2, \dots, A_s$  with each minimal polynomial  $m_{A_j}$  equals to a power  $(X - c_j)^{e_j}$  of a polynomial of degree 1 over  $k$ .*

□

**COROLLARY 4.6.** *A square matrix over a field  $k$  is semisimple if and only if it is diagonalizable (similar to a diagonal matrix) over  $\bar{k}$ .*

□

**COROLLARY 4.7.** *Let  $S_1, S_2$  be semisimple square matrices of the same size over  $k$ . If  $S_1$  and  $S_2$  commute, then both  $S_1 + S_2$  and  $S_1 S_2$  are also semisimple.*

**PROOF.** Using commutativity of  $S_1$  and  $S_2$ , we may easily see that  $S_1$  and  $S_2$  are simultaneously diagonalizable over  $\bar{k}$ . □

**COROLLARY 4.8.** *Let  $k$  be a field. Let  $C$  be a commutative subalgebra of  $M_n(k)$ . If  $C$  is generated by semisimple elements, then every element of  $C$  is also semisimple.*

□

On the other hand we have

**LEMMA 4.9.** *Let  $k$  be a field. Let  $C$  be a commutative subalgebra of  $M_n(k)$ . If  $C$  is generated by nilpotent elements, then every element of  $C$  is also nilpotent.*

**PROOF.** Easy. □

**COROLLARY 4.10.** *A Jordan-Chevalley decomposition (if there exists) of a square matrix  $A$  is unique.*

PROOF. Let

$$A = S + N = S' + N'$$

be two Jordan-Chevalley decompositions. Then  $S - S' = N' - N$  is a semisimple nilpotent element. Thus  $S - S' = N' - N = 0$ .  $\square$

PROOF. (of Proposition 4.2.) It now remains to prove that Jordan-Chevalley decomposition of a square matrix exists. By definition we may assume that  $k$  is algebraically closed. In view of Corollary 4.5, we may then assume that the minimal polynomial  $m_A$  of  $A$  is of the form  $(X - c)^e$  for some  $c \in k$  and  $e \in \mathbb{Z}_{>0}$ . Then

$$A = c + (A - c)$$

gives the required Jordan-Chevalley decomposition.  $\square$

DEFINITION 4.11. Let  $k$  be a field. For any square matrix  $x \in M_n(k)$ , we denote by  $x_s$  (respectively,  $x_n$ ) the semisimple (respectively, nilpotent) part of  $x$  in the Jordan-Chevalley decomposition of  $x$ .

LEMMA 4.12. *Let  $k$  be a field. Let  $x \in M_n(k)$  be a square matrix. then we have*

$$(\text{ad}(x))_s = \text{ad}(x_s), \quad (\text{ad}(x))_n = \text{ad}(x_n)$$

PROOF. Follows easily from the uniqueness of the Jordan-Chevalley decomposition.  $\square$

#### 4.2. $k$ -rationality.

PROPOSITION 4.13. *Let  $A$  be a square matrix over a field  $k$ . Let*

$$A = S + N$$

*be the Jordan-Chevalley decomposition of  $A$ .*

*Let  $m_A$  be the minimal polynomial of  $A$ . If all of the roots of  $m_A$  are separable over  $k$ , then  $S$  and  $N$  are defined over  $k$ . (That means, they are matrices over  $k$ ).*

PROOF. In view of Corollary , we may assume that  $m_A$  is of the form  $f^e$  for some irreducible polynomial  $f$  and a positive integer  $e$ . By assumption,  $f$  has only simple roots.

$$f(X) = \prod_{j=1}^d (X - c_j)$$



Let us define a polynomial  $\chi_s(X) \in \bar{k}[X]$  as follows.

$$\chi_s(X) = \frac{\prod_{\substack{1 \leq j \leq d \\ j \neq s}} (X - c_j)}{\prod_{\substack{1 \leq j \leq d \\ j \neq s}} (c_s - c_j)} \quad (s = 1, 2, 3, \dots, d)$$

These polynomials are designed to satisfy the following property.

$$\chi_s(c_j) = \begin{cases} 1 & \text{if } j = s \\ 0 & \text{if } j \neq s \end{cases}$$

Then we further define

$$\phi_s^{(e)}(X) = 1 - (1 - \chi_s^e)^e$$

and

$$\psi(X) = \sum_{s=1}^d c_s \phi_s^{(e)}(X).$$

It is fairly easy to see that

$$S = \psi(A)$$

holds.

The function  $\psi$  is symmetric with respect to roots  $\{c_s\}$  and thus  $\psi$  is a polynomial with coefficients in  $k$ . Thus  $S$  (hence also  $N$ ) is defined over  $k$ . □

The following example shows that the  $k$ -rationality of  $S$  does not necessarily hold when we drop off the assumption on  $A$ .

**EXAMPLE 4.14.** Let  $A \in M_p(\mathbb{F}_p(x))$  be a matrix of the following form.

$$A = \begin{pmatrix} 0 & 1 & & \\ & \ddots & \ddots & \\ & & \ddots & 1 \\ x & & & 0 \end{pmatrix}$$

Then the minimal polynomial of  $A$  is given by  $X^p - x$ . The Jordan-Chevalley decomposition of  $A$  is given by

$$A = x^{1/p} + (A - x^{1/p}).$$

Thus the decomposition is not defined over  $\mathbb{F}_p(x)$ .

## 5. GENERALITIES IN FINITE DIMENSIONAL LIE ALGEBRAS

## 5.1. Ideals of Lie algebras.

DEFINITION 5.1. For a linear subspace  $S, T$  of a Lie algebra  $L$ , let us define denote by  $[S, T]$  the following linear subspace of  $L$ .

$$(1) \quad [S, T] = (\text{linear span of } \{[x, y]; x \in S, y \in T\})$$

DEFINITION 5.2. Let  $L$  be a Lie algebra over a field  $k$ . A  $k$ -linear subspace  $\mathfrak{a}$  of  $L$  is said to be an **ideal** of  $L$  if

$$[x, y] \in \mathfrak{a} \quad (\forall x \in L, \forall y \in \mathfrak{a})$$

holds. This clearly is equivalent to saying that

$$[L, \mathfrak{a}] \subset \mathfrak{a}$$

holds.

PROPOSITION 5.3. *Let  $\mathfrak{a}$  be an ideal of a Lie algebra  $L$ . Then*

- (1)  $\mathfrak{a}$  is a sub  $L$ -module (sub representation) of  $L$ .
- (2)  $L/\mathfrak{a}$  carries a natural structure of a Lie algebra.

PROOF. As usual. □

## 5.2. Simple, semisimple, solvable, and nilpotent Lie algebras: definition.

DEFINITION 5.4. For a Lie algebra  $L$ , let us define the following ideals of  $L$ .

- (1)  $\text{Comm}(L) = [L, L]$ , and inductively,  

$$\text{Comm}^j(L) = \text{Comm}(\text{Comm}^{j-1}(L)).$$
- (2)  $\text{ad}(L)(L) = [L, L]$ , and inductively,  

$$\text{ad}^j(L)(L) = \text{ad}(\text{ad}^{j-1}(L)).$$

LEMMA 5.5. *For any Lie algebra  $L$  and for any positive integer  $j$ , we have*

$$\text{Comm}^j(L) \subset \text{ad}^j(L).$$

PROOF. Inductively, we have

$$\text{Comm}^j(L) = [\text{Comm}^{j-1}(L), \text{Comm}^{j-1}(L)] \subset [L, \text{ad}^{j-1}(L)] = \text{ad}^j(L).$$

□

DEFINITION 5.6. A Lie algebra  $L$  over a field  $k$  is said to be

- (1) **semisimple** if it has no abelian ideals.
- (2) **simple** if it has no non trivial ideals and  $\dim(L) > 1$ .

- (3) **solvable** if  $\text{Comm}^N(L) = 0$  for some  $N \in \mathbb{Z}_{>0}$ .
- (4) **nilpotent** if  $\text{ad}(L)^N(L) = 0$  for some  $N \in \mathbb{Z}_{>0}$ .

PROPOSITION 5.7. *We have the following implications.*

- (1) *Simple Lie algebras are semisimple.*
- (2) *Nilpotent Lie algebras are solvable.*

PROOF. (1) is Easy. (2) follows from Lemma 5.5. □

Semisimple algebras and solvable ones are “orthogonal”. For now we only mention the following

LEMMA 5.8. *Non zero solvable algebra  $L$  cannot be semisimple.*

PROOF. Let  $N_0$  be a positive integer such that

$$\text{Comm}_0^{N_0}(L) \neq 0, \quad \text{Comm}^{N_0+1}(L) = 0.$$

Then  $\text{Comm}^{N_0}(L)$  is a non-zero abelian ideal of  $L$ . □

### 5.3. The radicals of Lie algebras.

DEFINITION 5.9. A **radical** of a Lie algebra  $L$  is a maximal solvable ideal of  $L$ .

LEMMA 5.10. *Let  $\mathfrak{a}$  be an ideal of a Lie algebra  $L$ . If  $L/\mathfrak{a}$  and  $\mathfrak{a}$  are both solvable Lie algebras, then  $L$  is also solvable.*

PROOF. Since  $L/\mathfrak{a}$  is solvable, there exists a positive integer  $N_1$  such that

$$\text{Comm}^{N_1}(L/\mathfrak{a}) = 0.$$

Then we obviously have

$$\text{Comm}^{N_1}(L) \subset \mathfrak{a}.$$

On the other hand, since  $\mathfrak{a}$  is solvable, there exists a positive integer  $N_2$  such that

$$\text{Comm}^{N_2}(\mathfrak{a}) = 0.$$

We thus have

$$\text{Comm}^{N_1+N_2}(L) = \text{Comm}^{N_2}(\text{Comm}^{N_1}(L)) \subset \text{Comm}^{N_2}(\mathfrak{a}) = 0. \quad \square$$

LEMMA 5.11. *Every Lie subalgebras and quotients of solvable Lie algebras are solvable.*

PROOF. Obvious. □

LEMMA 5.12. *Let  $\mathfrak{a}, \mathfrak{b}$  be ideals of a Lie algebra  $L$ . If  $\mathfrak{a}, \mathfrak{b}$  are both solvable (as Lie algebras), then  $\mathfrak{a} + \mathfrak{b}$  is also solvable.*

PROOF.

$$\mathfrak{a} + \mathfrak{b}/\mathfrak{b} \cong \mathfrak{a}/\mathfrak{a} \cap \mathfrak{b}$$

□

PROPOSITION 5.13. *For a finite dimensional Lie algebra  $L$  over a field  $k$ , there exists a unique maximal solvable ideal of  $L$ . So we may call it **the radical** of  $L$ .*

PROOF. Let  $\mathfrak{a}_0$  be a solvable ideal of  $L$  which has the maximal dimension among solvable ideals. Then for any solvable ideal  $\mathfrak{b}$  of  $L$ ,  $\mathfrak{a}_0 + \mathfrak{b}$  is also solvable. Thus by the choice of  $\mathfrak{a}_0$  we see that

$$\mathfrak{a}_0 + \mathfrak{b} = \mathfrak{a}_0. \quad (\text{That is, } \mathfrak{a}_0 \supset \mathfrak{b}.)$$

Thus we see that  $\mathfrak{a}_0$  is the largest solvable ideal of  $L$ .

□

COROLLARY 5.14. *Let  $L$  be a finite dimensional Lie algebra over a field  $k$ . Let  $\mathfrak{r}$  be its radical. Then:*

- (1)  $L/\mathfrak{r}$  is semisimple.
- (2)  $L$  is semisimple if and only if  $\mathfrak{r} = 0$ .
- (3) A quotient  $L/\mathfrak{a}$  is semisimple if and only if  $\mathfrak{r} \subset \mathfrak{a}$ .

PROOF. (1) follows immediately from the definition and Lemma 5.8. (2) is also easy. (3):  $L/\mathfrak{a}$  contains

$$(\mathfrak{r} + \mathfrak{a})/\mathfrak{a} (\cong \mathfrak{r}/(\mathfrak{r} \cap \mathfrak{a}))$$

as a solvable ideal.

□

#### 5.4. Theorem of Engel.

LEMMA 5.15. *Let  $k$  be a commutative ring. Let  $x \in M_n(k)$  be a nilpotent matrix. Then  $\text{ad}(x)$  is also nilpotent.*

PROOF. Assume  $x^N = 0$ . We decompose  $\text{ad}(x)$  into left and right multiplication. Namely,

$$\text{ad}(x) = \lambda(x) - \rho(x)$$

Then  $\lambda(x)$  and  $\rho(x)$  commute with each other.

$$\text{ad}(x)^{2N-1} = \sum_{j=0}^{2N-1} \lambda(x)^j (-\rho(x))^{2N-1-j} = \sum_{j=0}^{2N-1} \lambda(x^j) \rho((-x)^{2N-1-j}) = 0.$$

□

THEOREM 5.16 (Engel). *Let  $V$  be a finite dimensional vector space over a field  $k$ . Let  $L$  be a Lie subalgebra of  $\mathfrak{gl}(V)$  such that each member of  $L$  is a nilpotent matrix. Then:*

- (1) If  $\dim(L) \geq 1$ , then  $L$  has an ideal of codimension 1.  
 (2) If  $\dim(V) \geq 1$ , then  $L$  has a simultaneous 0-eigen vector  $v$ .  
 (That is,  $x.v = 0 \quad (\forall x \in L), \quad v \neq 0$ .)

PROOF. If  $\dim(L) = 0$ , then there is nothing to do. We proceed by induction on  $\dim(L)$ . Let  $L_1$  be a maximal among

$$\{(\text{Lie subalgebras of } L \text{ which is not equal to } L)\}.$$

(The set above has a  $\{0\}$  as a member, so it is not empty.) In view of the lemma above, we see

$$\forall x \in L_1 \exists N \in \mathbb{Z}_{>0} (\text{ad}(x)^N(L) = 0).$$

We note that an vector space  $L/L_1$  admits adjoint actions by  $L_1$ . Thus

$$\forall x \in L_1 \exists N \in \mathbb{Z}_{>0} (\text{ad}(x)^N(L/L_1) = 0).$$

That means,  $\text{ad}_{L/L_1}(L_1) \subset \mathfrak{gl}(L/L_1)$  also satisfy the assumption of the theorem. By the induction hypothesis, we see that conclusion (2) is applicable to this case. Namely, there exists an element  $y_0 \in L \setminus L_1$  such that

$$\text{ad}(x)(y_0) (= [x, y_0]) \in L_1 \quad (\forall x \in L_1)$$

holds. Now a vector subset

$$L_2 = k.y_0 + L_1 (\supsetneq L_1)$$

of  $L$  is closed under Lie bracket and therefore it is a Lie subalgebra of  $L$ . By the maximality of  $L_1$ ,  $L_2$  should equal to  $L$ . It is then also easy to verify that  $L_1$  is an ideal of  $L (= L_2)$ . This proves (1).

To prove (2), we note that  $(L_1, V)$  satisfies the assumption of the theorem. So again by induction we see that  $L_1$  has a simultaneous 0-eigen vector. In other words,

$$V_1 := \bigcap_{x \in L_1} \{v \in V; x.v = 0\} \supsetneq \{0\}.$$

Let us then consider the action of  $y_0$ .

$$v \in V_1 \implies x.(y_0.v) = y_0.(x.v) + [x, y_0].v = 0 \implies y_0.v \in V_1$$

Thus  $V_1$  admits an action of  $y_0$ . Since  $y_0$  is nilpotent on  $V$  by the assumption, we see that  $y_0$  has at least one 0-eigen vector  $v_0 (\neq 0)$  in  $V_1$ . Then  $v_0$  surely is a simultaneous 0-eigen vector of  $L$ . □

**THEOREM 5.17 (Engel).** *Let  $V$  be a finite dimensional vector space over a field  $k$ . Let  $L$  be a Lie subalgebra of  $\mathfrak{gl}(V)$  such that each member*

of  $L$  is a nilpotent matrix. Then there exists a basis  $e_1, e_2, \dots, e_n$  of  $V$  such that

$$L \subset \mathfrak{n}_n(k) = \left\{ \begin{pmatrix} 0 & * & * & \dots & * \\ 0 & 0 & * & \dots & * \\ & & \ddots & * & * \\ 0 & 0 & \dots & 0 & * \\ 0 & 0 & \dots & 0 & 0 \end{pmatrix} \right\}$$

holds with respect to this basis. In particular,  $L$  is nilpotent.

PROOF. We apply the above theorem inductively to a vector space

$$L, L/(k.e_1), L/(k.e_1 + k.e_2) \dots$$

and obtain the desired basis  $\{e_i\}$ . Since the Lie algebra  $\mathfrak{n}_n(k)$  is nilpotent,  $L$  is also nilpotent. □

**5.5. Ideals of  $\mathfrak{gl}_n(k)$ .** Recall that  $\mathfrak{n}_n(k)$  denotes the Lie algebra of strictly upper triangular matrices. In this subsection we denote by  $e_{ij}$  the elementary matrices. (as we have done so without even mentioning...)

LEMMA 5.18. *Let  $k$  be a field of characteristic  $p$  (possibly 0). Let  $n \in \mathbb{Z}_{>1}$ . We assume that  $(p, n) \neq (2, 2)$ . Then we have*

$$\{x \in \mathfrak{gl}_n(k); [x, y] \in k.1_n \quad (\forall y \in \mathfrak{n}_n(k))\} = k.1_n + k.e_{1n}.$$

PROOF. Let us denote by  $L$  the left hand side of the lemma. Then we trivially have  $L \supset k.1_n$ . Furthermore, for all  $x \in \mathfrak{n}_n(k)$ , we see easily that

$$xe_{1n} = 0, \quad e_{1n}x = 0$$

holds. So we have

$$L \supset k.1_n + k.e_{1n}$$

Let us prove the opposite inclusion. We take an arbitrary element  $x \in L$ .

For any  $(i, j)$  satisfying  $i < j$ , we have  $e_{ij} \in \mathfrak{n}_n(k)$  and thus

$$[e_{i,j}, x] = c1_n \quad (\exists c \in k).$$

The rank of the left hand side is at most 2. So  $c$  must be equal to 0 when  $n \geq 3$ . Otherwise ( $n = 2$ ), we compare the trace of the both hand sides. The trace of the left hand side is clearly zero. The trace of a scalar matrix  $c.1_2$  is equal to  $2c$ . Thus  $c = 0$  by our assumption  $((p, n) \neq (2, 2))$ . In either case, we have  $[e_{i,j}, x] = 0$ . Then we compute

some of special cases. First, let us examine the case where  $i = 1, j \geq 2$ . Then

$$0 = [e_{1j}, x] = \sum_{st} [e_{1j}, x_{st}e_{st}] = \sum_t x_{jt}e_{1t} - \sum_s x_{s1}e_{sj}$$

By looking at  $(1, u)$  entry of the above equation, we conclude that equations in entries

$$\begin{aligned} \forall j \forall u ((j \geq 2, j \neq u) \implies x_{ju} = 0) \\ \forall j ((j \geq 2) \implies x_{jj} = x_{11}) \end{aligned}$$

hold. Similarly, by looking at the  $(u, n)$  entry of  $[e_{in}, x]$ , we conclude that equations

$$\forall i \forall u ((i \leq n-1, i \neq u) \implies x_{ui} = 0)$$

hold. Putting the equations all together, we conclude that  $x$  is in the right hand side of the lemma.  $\square$

As an application of the Engel's theorem, we prove the following proposition.

**PROPOSITION 5.19.** *Let  $k$  be a field of characteristic  $p$  (possibly 0). Let  $n$  be a positive integer. We assume that  $(p, n) \neq (2, 2)$ . Then each ideal  $I$  of  $\mathfrak{gl}_n(k)$  is equal to the one in the following list.*

- (1)  $0$ .
- (2)  $k \cdot 1_n$ .
- (3)  $\mathfrak{sl}_n(k)$ .
- (4)  $\mathfrak{gl}_n(k)$ .

**PROOF.** The case  $n = 1$  is trivial. So let us assume  $n \geq 2$ .

If  $I \subset k \cdot 1_n$ , then  $\dim_k(I) \leq 1$  and hence  $I = 0$  or  $I = k \cdot 1_n$ . Assume now  $I \not\subset k \cdot 1_n$ . Let us consider the Lie algebra  $\mathfrak{n}_n(k)$  of strictly upper triangular matrices. Then

$$(L, V) = (\mathfrak{n}_n(k), (I + k \cdot 1_n)/k \cdot 1_n)$$

satisfies the assumption of the Engel's theorem. So there exists a non-constant element  $x \in I$  such that

$$[x, y] \in k \cdot 1_n \quad (\forall y \in \mathfrak{n}_n(k)).$$

holds. By using the previous lemma, we see that  $x$  may be presented as

$$x = c_0 1_n + c_1 e_{1n} \quad (\exists c_0, c_1 \in k).$$

Since  $x$  is non-constant, we have  $c_1 \neq 0$ .

$$I \ni [e_{11}, x] = c_1 e_{1n}$$

Thus  $e_{1n}$  belongs to  $I$ . By changing the order of the base and repeating the above argument, we conclude that

$$\forall i \forall j ((i \neq j) \implies e_{ij} \in I.)$$

In addition we have

$$I \ni [e_{ij}, e_{ji}] = e_{ii} - e_{jj}$$

This clearly proves  $I \supset \mathfrak{sl}_n(k)$ . Since the codimension of  $\mathfrak{sl}_n(k)$  in  $\mathfrak{gl}_n(k)$  is 1, we have either  $I = \mathfrak{sl}_n(k)$  or  $I = \mathfrak{gl}_n(k)$ .  $\square$

For the sake of completeness, we deal with the case  $(p, n) = (2, 2)$ . In this case, situation is a bit different.

LEMMA 5.20. *Let  $k$  be a field of characteristic 2. Then:*

- (1) *Any two-dimensional vector subspace  $V$  of  $\mathfrak{sl}_2(k)$  with  $V \supset k \cdot 1_2$  is equal to a vector space  $L_{[b:c]}$  given by an element  $[b : c] \in \mathbb{P}^1(k)$  which is defined as*

$$L_{[b:c]} = k \cdot 1_n + k \cdot \begin{pmatrix} 0 & b \\ c & 0 \end{pmatrix}.$$

- (2) *For any element  $[b : c] \in \mathbb{P}^1(k)$ , the vector space  $L_{[b:c]}$  is an ideal of  $\mathfrak{gl}_2(k)$ .*
- (3) *In particular,  $\mathfrak{t}_2(k) = k \cdot 1_2 + \mathfrak{n}_2(k)$  is an ideal of  $\mathfrak{gl}_2(k)$ .*

PROOF. (1) There exists a traceless non constant matrix  $x \in V$  such that

$$I = k \cdot 1_n + k \cdot x$$

holds. By subtracting a constant matrix, one may easily replace  $x$  by a matrix with zero diagonals.

- (2) By a direct computation we see

$$\left[ \begin{pmatrix} x & y \\ z & w \end{pmatrix}, \begin{pmatrix} 0 & b \\ c & 0 \end{pmatrix} \right] = (x-w) \begin{pmatrix} 0 & b \\ c & 0 \end{pmatrix} + (bz-cy)1_2 \quad (x, y, z, w, b, c, \in k)$$

(Note that  $\text{char}(k) = 0$ .)

- (3)  $\mathfrak{t}_2(k) = L_{[1:0]}$ .  $\square$

PROPOSITION 5.21. *Let  $k$  be a field of characteristic 2. Then each ideal  $I$  of  $\mathfrak{gl}_2(k)$  is equal to the one in the following list.*

- (1)  $0$ .
- (2)  $k \cdot 1_n$ .
- (3) *A two-dimensional Lie algebra  $L_{[b:c]}$  defined as in the lemma above.*
- (4)  $\mathfrak{sl}_n(k)$ .



(5)  $\mathfrak{gl}_n(k)$ .

PROOF. We divide into several cases.

(i) Case where  $I \not\subset \mathfrak{sl}_2(k)$ . In this case there exists an element  $x \in I$  with  $\text{tr}(x) \neq 0$ . Putting

$$x = \begin{pmatrix} a & b \\ c & d \end{pmatrix},$$

we compute as follows

$$I \ni [e_{12}, x] = \left[ \begin{pmatrix} 0 & 1 \\ 0 & 0 \end{pmatrix}, \begin{pmatrix} a & b \\ c & d \end{pmatrix} \right] = \begin{pmatrix} c & a+d \\ 0 & c \end{pmatrix}.$$

(Note  $\text{char}(k) = 2$ .) Then we have

$$I \ni [e_{11}, [e_{12}, x]] = \left[ \begin{pmatrix} 1 & 0 \\ 0 & 0 \end{pmatrix}, \begin{pmatrix} c & a+d \\ 0 & c \end{pmatrix} \right] = \begin{pmatrix} 0 & a+d \\ 0 & 0 \end{pmatrix}$$

Thus we see that  $e_{12} \in I$ . In a same way (by changing the order of the base), we obtain,  $e_{21} \in I$ .

$$1_2 = [e_{21}, e_{12}] \in I.$$

Since  $\text{tr}(x) \neq 0$ , we see that  $\{1_2, e_{12}, e_{21}, x\}$  spans the  $\mathfrak{gl}_2(k)$ . thus  $I = \mathfrak{gl}_2(k)$  in this case.

(ii) Case where  $I \subset \mathfrak{sl}_2(k)$  and  $I \cap k.1_2 = 0$ . Let  $x$  be arbitrary element of  $I$  and put

$$x = \begin{pmatrix} a & b \\ c & d \end{pmatrix}.$$

Then by computing  $[e_{12}, x]$  as in the case (i) above, we know that  $c = 0$ . Similarly, we know  $b = 0$ . Since  $x$  is traceless,  $a = d$  also holds. So the only possibility in this case is  $I = 0$ .

(iii) Case where  $k.1_2 \subsetneq I \subsetneq \mathfrak{sl}_2(k)$ . By a dimension consideration, we see that  $\dim I = 2$ . Then we use the above lemma.

(iv) The case  $I = k.1_2$  or  $I = \mathfrak{sl}_2(k)$ . Excellent. There is nothing to in this case.

□

### 5.6. Ideals of $\mathfrak{sl}_n(k)$ .

PROPOSITION 5.22. *Let  $k$  be a field of characteristic  $p$  (possibly 0.) Let  $n$  be a positive integer.*

- (1) *If  $p \nmid n$ , then  $\mathfrak{sl}_n(k)$  is a simple Lie algebra.*
- (2) *If  $p \mid n$ , then  $\mathfrak{sl}_n(k)$  has a unique nontrivial ideal  $k.1_n$ .*

PROOF. Let  $I$  be an ideal of  $\mathfrak{sl}_n(k)$ . By taking trace we see immediately that

$$1_n \in \mathfrak{sl}_n(k) \iff p|n.$$

Thus if  $I \subset k \cdot 1_n$ , then The only nontrivial possibility is that  $p|n$  and  $I = k \cdot 1_n$ .

Assume now that  $I \not\subset k \cdot 1_n$ . Then by an argument similar to that in Proposition 5.19, we see that

$$x = c_0 1_n + c_1 e_{1n} \in I, \quad (\exists c_0, c_1 \in k, c_1 \neq 0)$$

holds.

(1) If  $p \nmid n$ , then by taking trace we see that  $e_{1n} \in I$ . By permuting the basis, we see that  $e_{ij} \in I$  whenever  $i \neq j$ . Thus  $I = \mathfrak{sl}_n(k)$  in this case.

(2) If  $p|n$ , then by assumption on  $(p, n)$  we have  $n \geq 3$ . Thus

$$I \ni [e_{21}, x] = c_1 e_{2n} \quad \therefore e_{1n} = [e_{12}, e_{2n}] \in I.$$

So in this case also we see that  $I = \mathfrak{sl}_n(k)$ . □

PROPOSITION 5.23. *Let  $k$  be a field of characteristic 2. Then any two dimensional Lie algebra  $L_{[b:c]}$  as in Lemma 5.20 is an ideal of  $\mathfrak{sl}_2(k)$ . Thus each ideal  $I$  of  $\mathfrak{sl}_2(k)$  is equal to the one in the following list.*

- (1) 0.
- (2)  $k \cdot 1_n$ .
- (3)  $L_{[b:c]}$
- (4)  $\mathfrak{sl}_n(k)$

PROOF. Easy exercise. □

### 5.7. Invariant bilinear forms and Killing forms.

DEFINITION 5.24. A symmetric bilinear form  $B : L \times L \rightarrow k$  of a Lie algebra over a field  $k$  is said to be **invariant** if it satisfies

$$B([Y, X], Z) + B(X, [Y, Z]) = 0 \quad (\forall X, Y, Z \in L)$$

(which means that “the Lie derivative of  $B$  is zero”), or, equivalently,

$$B([X, Y], Z) = B(X, [Y, Z]) \quad (\forall X, Y, Z \in L)$$

(which means that  $B$  is “balanced”).

LEMMA 5.25. *Let  $L$  be a Lie algebra over a field  $k$ . Let  $B$  be an invariant bilinear form on  $L$ . Then for any ideal  $I$  of  $L$ ,*

$$L^\perp = \{x \in L; B(x, y) = 0(\forall y \in l)\}$$

*is an ideal of  $L$ .*

PROOF. Easy. □

Note: We need to be a bit careful when we use the notation  $\bullet^\perp$ . It is safer to clarify the “container” ( $L$ ) and bilinear form  $\rho$ . So the lemma above we should have written  $L^{\perp\rho,L}$  (eek) in stead of  $L^\perp$ .

DEFINITION 5.26. Let  $(\rho, V)$  be a finite dimensional representation of a Lie algebra  $L$  over a field  $k$ . Then the **Killing form with respect to  $(\rho, V)$**  is a bilinear form on  $L$  defined by

$$\text{Tr}_{\rho,V}(XY) = \text{tr}_V(\rho(X)\rho(Y)).$$

The **ordinary(usual) Killing form**  $\kappa_L$  of  $L$  is a bilinear form on  $L$  defined as the Killing form of the adjoint representation. That is,

$$\kappa_L(X, Y) = \text{Tr}_{\text{ad},L}(XY) = \text{tr}_{\text{ad},L}(\text{ad}(X) \text{ad}(Y)).$$

It is easy to verify that the Killing forms defined as above are invariant.

5.7.1. *functoriality of Killing forms.*

LEMMA 5.27. *Let  $L$  be a Lie algebra over a field  $k$ . Then the followings are true.*

- (1) *Let  $V$  be a finite dimensional representation of  $L$ . Let  $W$  be a subrepresentation of  $V$ . Then we have*

$$\text{Tr}_V(xy) = \text{Tr}_W(xy) + \text{Tr}_{V/W}(xy).$$

- (2) *Let  $I$  be an ideal of  $L$ . Assume  $L$  is finite dimensional. Then we have*

$$\kappa_L(x, y) = \text{Tr}_{\text{ad},L}(xy) = \text{Tr}_{\text{ad},I}(xy) + \text{Tr}_{\text{ad},L/I}(xy) = \text{Tr}_{\text{ad},I}(xy) + \kappa_{L/I}(\bar{x}, \bar{y})$$

*(where  $\bar{\bullet}$  denotes the class of  $\bullet$  in  $L/I$ .) In particular, for any  $x, y \in I$ , we have*

$$\kappa_L(x, y) = \kappa_I(x, y)$$

PROOF. (2): We choose a basis  $B = B_1 \amalg B_2$  of  $L$  such that  $B_2$  forms a basis of  $I$ . Then  $\bar{B}_1$  forms a basis of  $L/I$ . Under the basis  $B$ ,  $\text{ad}(x)$  may be represented by a matrix

$$\text{ad}_L(x) = \begin{pmatrix} \text{ad}_{L/I}(\bar{x}) & * \\ 0 & \text{ad}_I x \end{pmatrix}.$$

We obtain the result easily from this.

(1): may be proved in a same manner. □

### 5.8. Theorem of Iwasawa.

**THEOREM 5.28.** *Let  $L$  be a Lie algebra over a field  $k$  of characteristic  $p \neq 0$ . Then  $L$  has a finite dimensional faithful representation. More precisely, there exists a two-sided ideal  $I$  of the universal enveloping algebra  $U(L)$  such that  $L$  acts faithfully on  $U(L)/I$ .*

Before proving the above theorem, we first prove the next lemma.

**LEMMA 5.29.** *Under the hypothesis of the theorem, for any  $x \in L$ , there exists a monic non constant polynomial  $f_x(X) \in k[X]$  such that*

$$f_x(x) \in Z(U(L))$$

*holds.*

**PROOF.** Let us put  $s = \dim(L)$ . The linear transformation  $\text{ad}(x)$  on  $L$  is represented by a matrix of size  $s$  and has therefore its minimal polynomial  $m_x$ : Namely,  $m_x$  is a monic polynomial of degree no more than  $s$  such that

$$m_x(\text{ad}(x)) = 0$$

holds. Let us divide  $X, X^p, X^{p^2}, \dots, X^{p^{s+1}}$  by  $m_x(X)$ .

$$X^{p^j} = m_x(X)q_j(X) + r_j(X) \quad (\deg(r_j) < s) \quad (j = 1, 2, 3, \dots, s+1)$$

Then  $s+1$  polynomials  $\{r_j(X)\}_{j=1}^{s+1}$  of degree  $\leq s-1$  should be linearly dependent. That means, there exists a non trivial vector  $(c_j) \in k^{s+1}$  such that

$$\sum_j c_j X^{p^j} \in m_x(X)k[X]$$

holds. Then we have

$$\sum_j c_j (\text{ad}(x))^{p^j} = 0.$$

Thus we conclude

$$\text{ad}\left(\sum_j c_j x^{p^j}\right) = 0.$$

By dividing  $\sum_j c_j X^{p^j}$  by leading coefficient, we obtain the required polynomial  $f_x$ . □

**PROOF.** of the Theorem Let  $\{e_1, e_2, \dots, e_s\}$  be a basis of  $L$ . Then by the above lemma we know that there exists a set of monic non constant polynomials  $\{f_1, f_2, \dots, f_s\}$  such that each  $h_j = f_j(e_j)$  belongs to the

center of  $U(L)$ . Let us put  $d_j = \deg(f_j)$ . Then using PWB theorem we may easily see that

$$\left\{ \begin{array}{l} c_1, c_2, c_3, \dots, c_s \in \mathbb{N}, \\ h_1^{c_1} h_2^{c_2} h_3^{c_3} \dots h_s^{c_s} e_1^{l_1} e_2^{l_2} e_3^{l_3} \dots e_s^{l_s}; l_1, l_2, l_3, \dots, l_s \in \mathbb{N}, \\ l_j < d_j (\forall j) \end{array} \right\}$$

forms a basis of  $U(L)$ . Let us now put

$$I = U(L)(h_1, \dots, h_s)$$

Then  $A = U(L)/I$  is a finite dimensional vector space with the base

$$\left\{ e_1^{l_1} e_2^{l_2} e_3^{l_3} \dots e_s^{l_s}; l_1, l_2, l_3, \dots, l_s \in \mathbb{N}, l_j < d_j (\forall j) \right\}$$

The representation  $\rho_A$  of  $L$  on  $A$  is faithful. Indeed, for any  $x \in L$ , we have

$$\rho_A(x) = 0 \implies x.1 = 0 \text{ in } A \implies x \in I \implies x = 0.$$

□

DEFINITION 5.30. Let  $L$  be a Lie algebra.

- (1) A representation  $V$  of  $L$  is called **completely reducible** if it is a direct sum of reducible sub representations.
- (2)  $L$  is called **completely reducible** if every representation of  $L$  is completely reducible.

The following remark is (at least) in the Book of Bourbaki.

PROPOSITION 5.31. *Let  $L$  be a non zero finite dimensional Lie algebra over a field  $k$  of characteristic  $p \neq 0$ . Then  $L$  can never be completely reducible.*

PROOF. Let us follow the proof of the theorem of Iwasawa. By taking  $f_1^2$  instead of  $f_1$  in the proof, we obtain a representation  $A_1 = U(L)/I$  with a non trivial central nilpotent  $z = f_1(e_1)$ . Then we see that  $zA_1$  cannot have a direct complementary  $L$ -module  $X$ . For if it existed, then  $X$  should necessarily a left ideal of  $A_1$ . On the other hand, by decomposing  $1 \in A_1$  we obtain

$$1 = x + za \quad (\exists x \in X \exists a \in A_1).$$

Then  $x = 1 - za$  has an inverse  $(1 + za)$ . This implies that

$$X \supset A_1 x \supset A_1.$$

which is a contradiction. □

**5.9. Cartan's criterion for solvability (Ccs).** Cartan's criterion relates several properties (semi simplicity, solvability) of Lie algebras with properties of their invariant bilinear forms.

To prove it we need some study on bilinear forms.

**DEFINITION 5.32** (in this subsection only). For any free module  $R^t$  over a ring  $R$ , We denote by  $\langle \rangle_R$  the "usual inner product". That is,

$$\langle v, w \rangle_R = \sum_i v_i w_i.$$

The first thing we do is to observe the property of this inner product when the base ring  $R$  is a "real field". (Since we only need it for the case  $R = \mathbb{Q}$ , we omit the definition of a real field and describe the following lemma only when  $R = \mathbb{Q}$ .)

**LEMMA 5.33.** *Let  $W = \mathbb{Q}^t$  be a vector space. Let  $b_1, b_2, \dots, b_s \in W$ . Let  $B$  be a  $s \times s$  matrix defined by*

$$B = (\langle b_i, b_j \rangle_{\mathbb{Q}}).$$

*Then the determinant of  $B$  is equal to 0 if and only if  $\{b_j\}_{j=1}^s$  are linearly dependent over  $\mathbb{Q}$ .*

**PROOF.** Assume  $\{b_j\}_{j=1}^s$  are linearly dependent over  $\mathbb{Q}$ . Then there exists a non trivial vector  $(c_1, c_2, \dots, c_s) \in \mathbb{Q}^s$  such that

$$\sum_{j=1}^s c_j b_j = 0$$

holds. Thus

$$(c_1, c_2, \dots, c_s)B = 0$$

So  $B$  is a degenerate matrix which implies that  $\det(B) = 0$ .

Let us now prove the opposite implication. Assume  $\det(B) = 0$ . Then there exists a non trivial vector  $(c_1, c_2, \dots, c_s)$  such that

$$(c_1, c_2, \dots, c_s)B = 0$$

holds. Let us put

$$v = \sum_j c_j b_j.$$

Then we see that  $\langle v, v \rangle_{\mathbb{Q}} = 0$  and hence  $v = 0$ . (Note that for this implication we have used the fact that  $\mathbb{Q}$  is a "real field".) Thus  $\{b_j\}$  are linearly dependent over  $\mathbb{Q}$ .  $\square$

The next task is to compare  $\mathbb{Q}$  with other field.

DEFINITION 5.34. For any subset  $S$  of a  $\mathbb{Z}$ -module  $W_{\mathbb{Z}}$ , Let us put

$$\text{MG}_S = \max\{|\det(\langle b_i, b_j \rangle_{\mathbb{Z}})| \mid b_1, \dots, b_l \in S\}.$$

(“The maximum modulus of Gram determinants”.) We denote by  $S_k$  the subset of  $W_k = W_{\mathbb{Z}} \otimes_{\mathbb{Z}} k$  defined by

$$S_k = \{x \otimes 1 \in W \otimes_{\mathbb{Z}} k; x \in S\}.$$

LEMMA 5.35. Let  $S$  be a finite subset of a free module

$$W_{\mathbb{Z}} = \mathbb{Z}^t = \{(v_1, v_2, \dots, v_t); v_j \in \mathbb{Z}(\forall j)\}$$

over  $\mathbb{Z}$ . Let  $k$  be a field of characteristic  $p$ . We assume either  $p = 0$  or  $p > \text{MG}_S$  holds. Then we have

$$(a \in W_k, (a^{\perp} \cap S_k)^{\perp\perp} \ni a) \implies a = 0$$

PROOF. Assume  $a \neq 0$ . Since the inner product  $\langle \bullet, \bullet \rangle_k$  is non degenerate on  $W_k$ , we see that  $(a^{\perp} \cap S_k)^{\perp\perp}$  is equal to the  $k$ -vector space spanned by  $(a^{\perp} \cap S_k)$ . Thus there exists a set of linearly independent vectors  $\{b_j\} \subset S_k$  so that we may write down  $a$  as

$$a = \sum_i a_i b_i \quad (a_i \in k).$$

Then by the assumption on  $a$ , we see that

$$\sum a_i \langle b_i, b_j \rangle_k = 0 \quad (\forall j)$$

Thus

$$\det(\langle b_i, b_j \rangle_k) = 0$$

which is equivalent to

$$(1) \quad p \mid \det(\langle b_i, b_j \rangle_{\mathbb{Z}}).$$

Note on the other hand that  $b_j$  are linearly independent over  $\mathbb{Z}$ . Thus

$$\det(\langle b_i, b_j \rangle_{\mathbb{Z}}) \neq 0$$

By the definition of  $\text{MG}_S$ , we see that

$$(2) \quad 0 < |\det(\langle b_i, b_j \rangle_{\mathbb{Z}})| \leq \text{MG}_S$$

which contradicts to the condition (1).  $\square$

DEFINITION 5.36. For any positive integer  $n$ , and for any ring  $k$ , we denote by  $\text{Diag}_n(k)$  the set of diagonal matrices in  $M_n(k)$ . For any vector  $a = (a_i) \in k^n$ , we denote by  $\text{diag}(a)$  the diagonal matrix  $\text{diag}(a) = \text{diag}(a_1, \dots, a_n)$ . Note that for any ring  $k$ , the restriction of the Killing form of  $\mathfrak{gl}_n$  coincides with the “usual” inner product with this identification. That is,

$$\text{tr}(\text{diag}(a) \text{diag}(b)) = \langle a, b \rangle_k.$$

We define the following subset of  $\text{Diag}_n(\mathbb{Z})$ .

$$S^{[n]} = \{(\text{diag}((e_i - e_j) - (e_m - e_l)); i, j, m, l \in \{1, 2, 3, \dots, n\})\}.$$

(where the vectors  $\{e_i\}_{i=1}^n$  are elementary vectors.) We note that an obvious estimate

$$\text{MG}_{S^{[n]}} \leq 4^n n!$$

holds.

**LEMMA 5.37.** *Let  $n$  be a positive integer. Let  $a = (a_1, \dots, a_n), b = (b_1, \dots, b_n) \in k^n$ .*

*If  $b$  satisfies  $b \in (a^\perp \cap S^{[n]})^\perp$ , then there exist polynomials  $f, g \in k[X]$  such that*

$$f(\text{diag}(a)) = \text{diag}(b), \quad g(\text{ad}(\text{diag}(a))) = \text{ad}(\text{diag}(b))$$

*holds.*

**PROOF.** Let us denote by  $\epsilon_{ijlm}$  the vector

$$\epsilon_{ijlm} = (e_i - e_j) - (e_l - e_m).$$

We first find a map  $f_0$  from  $\Lambda_a = \{a_i; i = 1, 2, \dots, n\}$  to  $\Lambda_b = \{b_i; i = 1, 2, \dots, n\}$  such that

$$f_0(a_i) = b_i.$$

Such a thing exists (is “well defined”) if and only if

$$\forall i \forall j (a_i = a_j \implies b_i = b_j)$$

holds. This condition is equivalent to the condition

$$\forall i \forall j (a \perp \epsilon_{ij11} \implies b \perp \epsilon_{ij11})$$

which holds by the assumption on  $b$ . Thus we see that  $f_0$  exists. On the other hand, by using Lagrange interpolation formula we see that there exists a polynomial  $f \in k[X]$  such that  $f|_{\Lambda_a} = f_0$ . Then we have

$$f(\text{diag}(a)) = \text{diag}(b).$$

The adjoint action of a diagonal matrix  $\text{diag}(a)$  is represented by a diagonal matrix  $(a_i - a_j)_{i,j}$ . Thus an argument similar to the one above proves the existence of  $g$ . □

**PROPOSITION 5.38 (Cartan).** *Let  $V$  be an  $n$  dimensional vector space over a field  $k$  of characteristic  $p$ . We assume that either  $p = 0$  or  $p > \text{MG}_{S^{[n]}}$  holds. Let  $L$  be a Lie subalgebra of  $\mathfrak{gl}(V)$ . If the Killing form of  $L$  with respect to  $V$  is identically equal to 0, then  $L$  is solvable.*



PROOF. We may assume that  $k$  is algebraically closed. Let us take an element  $x \in [L, L]$ . Then we have

$$(\text{ad } x_s)(L) = (\text{ad } x)_s(L) \subset L$$

Let us now diagonalize  $x_s$  and write  $x_s = \text{diag}(a)$ . Let us take arbitrary  $b \in (a^\perp \cap S^{[n]})^\perp$ . By the lemma above we see that there exist polynomials  $f, g \in k[X]$  such that

$$f(\text{diag}(a)) = \text{diag}(b), \quad g(\text{ad}(\text{diag}(a))) = \text{ad}(\text{diag}(b))$$

holds. For any  $w = \sum_l [y_l z_l] \in [L, L]$ , we have:

$$\begin{aligned} \text{tr}(\text{diag}(b)[\sum_l y_l z_l]) &= \sum_l \text{tr}([\text{diag}(b), y_l] z_l) \\ &= \sum_l \text{tr}(\text{ad}(\text{diag}(b)).y_l z_l) = \sum_l \text{tr}(g(\text{ad}(\text{diag}(a))).y_l z_l) \\ &= \sum_l \text{tr}(g(\text{ad}(x_s)).y_l z_l) \in \sum_l \text{tr}(LL) \\ &= 0 \end{aligned}$$

That means,  $\text{tr}(\text{diag}(b)w) = 0$ . In particular, we have

$$\text{tr}(\text{diag}(b)x) = 0.$$

Since  $\text{diag}(b) = f(\text{diag}(a)) = f(x_s)$  is a polynomial in  $x$ , it commutes with  $x_s$  and with  $x_n$ . thus

$$\text{diag}(b)x = (\text{diag}(b)x_s) + (\text{diag}(b)x_n)$$

gives the Jordan-Chevalley decomposition of  $\text{diag}(b)x$ . Therefore,

$$0 = \text{tr}(\text{diag}(b)x) = \text{tr}(\text{diag}(b)x_s) = \text{tr}(\text{diag}(b) \text{diag}(a)) = \langle b, a \rangle_k$$

thus  $b \perp a$ .

To sum up, we have shown

$$(a^\perp \cap S)^\perp \ni b \implies b \in a^\perp.$$

In other words,

$$(a^\perp \cap S)^{\perp\perp} \ni a$$

which is equivalent to saying that  $a$  is a linear combination of elements in  $(a^\perp \cap S)$ .

In view of Lemma 5.35, we see that  $a = 0$ . So  $x = x_s + x_n = x_n$  is a nilpotent element.

By the theorem of Engel, we conclude that  $[L, L]$  is nilpotent. Thus  $L$  is solvable (since we have shown that  $L/[L, L]$  and  $[L, L]$  are solvable).  $\square$

DEFINITION 5.39. We say that the Cartan's criterion for solvability (Ccs) holds for a linear Lie algebra  $L \subset (\mathfrak{gl}_n(k))$  over a field  $k$  if it satisfies the following condition.

(Ccs) If the Killing form on  $L$  associated to  $k^n$  is identically zero,  $L$  is solvable.

Let  $n$  be a positive integer. We denote by  $\text{Ccs}(n)$  the set of  $p$  such that Ccs holds for any Lie algebra  $L$  of dimension less than or equal to  $n$  for any field  $k$  of characteristic  $p$ .

$$\text{Ccs}(n) = \{p; \text{Ccs holds for any } (L \subset \mathfrak{gl}_n(k)) \text{ provided } \text{char}(k) = p, \}$$

COROLLARY 5.40 (of Proposition). *Let  $n$  be a positive integer. Then:*

- (1)  $0 \in \text{Ccs}(n)$
- (2) *For any prime  $p$  which is larger than  $\text{MG}_{S^{[n]}}$ , we have  $p \in \text{Ccs}(n)$ .*
- (3) *In particular, for any prime  $p$  which is larger than  $4^n n!$ , we have  $p \in \text{Ccs}(n)$ .*

Note:

The estimate given in the above corollary is presumably far from the best one.

PROPOSITION 5.41. *Let  $n$  be a positive integer. Let  $k$  be a field of characteristic  $p \in \text{Ccs}(n)$ . Let  $L$  be a Lie algebra over  $k$  whose dimension is less than or equal to  $n$ . If the usual Killing form  $\kappa = \text{Tr}_{\text{ad},L}$  of  $L$  is identically equal to zero, then  $L$  is solvable. In particular, if  $p > 4^n n!$  or  $p = 0$ , then  $L$  is solvable if its usual Killing form is identically equal to zero.*

PROOF. Apply the definition to

$$L/(\text{center of } L) \hookrightarrow \mathfrak{gl}_n(L).$$

□

### 5.10. Cartan's criterion for semisimplicity.

DEFINITION 5.42. We call a Lie algebra  $L$  over  $k$  **non degenerate** if the Killing form  $\kappa_L$  of  $L$  is non degenerate.

LEMMA 5.43. *Every non degenerate Lie algebra  $L$  over a field  $k$  is semisimple.*

PROOF. Assume that there exists a non trivial abelian ideal  $A$  of  $L$ . Let  $y_0$  be a non zero element of  $A$ . Then for any  $x \in L$ ,  $z = \text{ad}(x) \text{ad}(y_0)$  is nilpotent. Indeed,

$$z(L) = \text{ad}(x) \text{ad}(y_0)(L) = \text{ad}(x)([y_0, L]) \subset \text{ad}(x)(A) \subset A,$$

$$z^2(L) = \text{ad}(x) \text{ad}(y_0)(z(L)) \subset \text{ad}(x) \text{ad}(y_0)(A) = \text{ad}(x)[y_0, A] = 0.$$

Thus  $\kappa(x, y_0) = \text{tr}(\text{ad}(x) \text{ad}(y_0)) = 0$  for any  $x \in L$ . That means,  $y_0 \in L^\perp$ . This is contrary to the assumption on  $L$ .  $\square$

**PROPOSITION 5.44.** *Let  $n$  be a positive integer. Let  $k$  be a field of characteristic  $p \in \text{Ccs}(n)$ . Let  $L$  be a Lie algebra of dimension  $n$ . Then the following conditions are equivalent:*

- (1)  $L$  is semisimple.
- (2)  $L$  is non degenerate.
- (3)  $L$  is a direct sum of simple ideals.

**PROOF.** ((1)  $\implies$  (2)): Assume  $L$  is semisimple. Let us take an ideal  $I = L^\perp$  of  $L$ . Then the Killing form on  $I$  is identically equal to zero. since  $\dim(I) \leq n$ ,  $I$  is a solvable algebra. Since  $L$  is semisimple, this implies  $I = 0$ .

((2)  $\implies$  (1)): holds (regardless of the base field) in view of the previous lemma.

((3)  $\implies$  (2)): We see that simple algebras are non degenerate in view of the argument above. Thus  $L$  is also non degenerate.

((2)  $\implies$  (3)): Let  $H$  be a nontrivial ideal of  $L$ . Then  $H \cap H^\perp$  is an abelian ideal of  $L$ . Indeed, for any  $y, z \in H \cap H^\perp$  and for any  $x \in L$ , we have

$$\kappa(x, [y, z]) = \kappa([x, y], z) \in \kappa(H, H^\perp) = 0$$

So that  $[y, x] \in L^\perp = 0$ . On the other hand, by the previous lemma we see that  $H \cap H^\perp$  is semisimple and so we have  $H \cap H^\perp = 0$ . Accordingly we have  $L = H \oplus H^\perp$ .  $\square$

### 5.11. examples.

**EXAMPLE 5.45.** Let  $k$  be a field of characteristic  $p$  (possibly 0).

$$\mathfrak{gl}_n(k)$$

is a Lie algebra with the Killing form

$$\begin{aligned} \kappa_{\mathfrak{gl}_n(k)}(x, y) &= \text{tr}((\lambda(x) - \rho(x))(\lambda(y) - \rho(y))) \\ &= \text{tr}(\lambda(xy)) + \text{tr}(\rho(xy)) - \text{tr}(\lambda(x)\rho(y)) - \text{tr}(\lambda(y)\rho(x)) \\ &= 2n \text{tr}(xy) - 2 \text{tr}(x) \text{tr}(y). \end{aligned}$$

$\mathfrak{sl}_n(k)$  is an ideal of  $\mathfrak{gl}_n(k)$  and so its Killing form is given by

$$\kappa_{\mathfrak{sl}_n(k)}(x, y) = 2n \text{tr}_k(xy).$$

If  $p \nmid 2n$ , then the Killing form is easily seen to be non-degenerate. so  $\mathfrak{sl}_n(k)$  is a non-degenerated Lie algebra in this case. In this way we see

that it is a semisimple Lie algebra. The Lie algebra is actually simple as we have shown in Proposition 5.22.

EXAMPLE 5.46. Let  $p$  be an odd prime. Let  $k$  be a field of characteristic  $p$ . Then we have shown in Proposition 5.19 that the only nontrivial ideals of  $\mathfrak{gl}_p(k)$  are  $\mathfrak{sl}_p(k)$  and  $k \cdot 1_p$ . So we see that

$$L = \mathfrak{gl}_p(k)/k \cdot 1_p$$

is a semisimple Lie algebra (as it has no proper abelian ideals). It has a unique nontrivial ideal

$$M = \mathfrak{sl}_p(k)/k \cdot 1_p.$$

Thus  $L$  cannot be a direct sum of simple Lie algebras.

### 5.12. Weyl's theorem on complete reducibility.

DEFINITION 5.47. Let  $L$  be a finite dimensional Lie algebra. Let  $B$  be a non-degenerate invariant bilinear form on  $L$ . Then we define the **Casimir element**  $C_B \in U(L)$  with respect to  $B$  by

$$C_B = \sum_i x_i x^{(i)}$$

where  $\{x_i\}$  is a basis of  $L$ , and  $\{x^{(i)}\}$  is the dual basis of the basis  $\{x_i\}$  with respect to  $B$ .

PROPOSITION 5.48. *Under the same assumption of the definition above, we have the following facts.*

- (1) *The Casimir operator  $C_B$  is independent of the choice of the basis  $\{x_i\}$  of  $L$ .*
- (2)  *$C_B$  commutes with  $L$ . So it is in the center of  $U(L)$ .*

PROOF. (1): easy exercise in linear algebra.

(2): For any  $a \in L$ , let us write the adjoint action of  $a$  on  $L$  by using the basis  $\{x_i\}$ . Namely,

$$[a, x_i] = \sum_j c_i^{(j)}(a) x_j \quad (c_i^{(j)}(a) \in k).$$

Then the constants  $\{c_i^{(j)}(a)\}$  (“structure constants”) are expressed in terms of  $B$  as follows.

$$B(x^{(l)}, [a, x_i]) = \sum_j c_i^{(j)}(a) B(x^{(l)}, x_j) = c_i^{(l)}(a)$$

We note that from the invariance of  $B$ , we have

$$B([x^{(l)}, a], x_i) = c_i^{(l)}(a),$$

so that we have a dual expression

$$[x^{(l)}, a] = \sum_i c_i^{(l)}(a)x^{(i)}.$$

Then we compute as follows.

$$\begin{aligned} [a, C_B] &= \sum_i [a, x_i]x^{(i)} + \sum_i x_i[a, x^{(i)}] \\ &= \sum_i \sum_j c_i^{(j)}(a)x_jx^{(i)} - \sum_i \sum_j c_j^{(i)}(a)x_ix^{(j)} = 0. \end{aligned}$$

□

DEFINITION 5.49. Let  $L$  be a finite dimensional Lie algebra. Let  $V$  be a finite dimensional  $L$ -module.  $\text{Tr}_V(\bullet\bullet)$  with respect to  $V$ . We assume that the Killing form  $\text{Tr}_V(\bullet\bullet)$  with respect to  $V$  is non degenerate. Then we define the **Casimir element** with respect to  $V$  by

$$C_V = C_{\text{Tr}_V}.$$

We see immediately from the definition that

$$\text{Tr}_V(C_V) = \dim(V)$$

holds in general.

LEMMA 5.50. *Let  $k$  be a field of characteristic  $p$ . Let  $L$  be a  $n$ -dimensional semisimple Lie algebra over a field  $k$ . Let  $V$  be a  $m$ -dimensional  $L$ -module. Let  $I$  be the kernel of the representation  $\rho_V$  associated to  $V$ . We assume  $p \in \text{Ccs}(n) \cap \text{Ccs}(m)$ . Then the Killing form  $\text{Tr}_V$  on  $L/I$  is non degenerate.*

PROOF.  $L$  is semisimple and  $p \in \text{Ccs}(n)$  so  $L$  is non degenerate.  $L/I$  is also non-degenerate so  $L/I$  is semisimple. We may thus assume  $I = 0$ . An ideal

$$J = L^{\perp_{\text{Tr}_V}}$$

of  $L$  is a solvable ideal. Since  $L$  is semisimple and  $p \in \text{Ccs}(m)$ , we have by  $J = 0$ . That means,  $\text{Tr}_V$  is non-degenerate on  $L$ .

□

LEMMA 5.51. *Let  $k$  be a field of characteristic  $p$  (which may be 0). Let  $(L, W \subset V)$  be a triple which satisfies the following conditions.*

- (1)  $L$  is a finite dimensional semisimple Lie algebra over  $k$ .
- (2)  $V$  is a finite dimensional  $L$ -module.
- (3)  $W$  is an  $L$ -submodule of  $V$  of codimension 1.
- (4)  $p \in \text{Ccs}(\dim(L)) \cap \text{Ccs}(\dim(V))$ .
- (5)  $p > \dim(V)$  or  $p = 0$ .

Then the exact sequence

$$0 \rightarrow W \rightarrow V \rightarrow V/W \rightarrow 0$$

splits. In other words, there exists a 1-dimensional  $L$ -submodule  $X$  of  $V$  which is complementary to  $W$ .

PROOF. Since the question of existence of  $X$  is described in terms of existence of a solution of a set of linear equations, we may assume that  $k$  is algebraically closed. Let us denote by  $\rho_V$  the representation of  $L$  associated to  $V$ . Then by replacing  $L$  by  $L/\ker(\rho_V)$  if necessary, we may assume that the representation  $\rho_V$  is faithful.

Note that since  $L$  is semisimple, it acts on  $V/W$  trivially.

Let us first treat the case where  $W$  is irreducible. Let  $c = C_V$  be a Casimir element with respect to  $V$ . Since  $L$  acts on  $V/W$  trivially,  $c|_{V/W}$  is equal to zero. Thus

$$\dim(V) = \operatorname{tr}_V(c) = \operatorname{tr}_W(c|_W) + \operatorname{tr}_{V/W}(c|_{V/W}) = \operatorname{tr}_W(c|_W).$$

In particular,  $c|_W$  is not equal to zero. On the other hand, by Schur's lemma,  $c|_W$  is equal to a scalar  $\lambda \in k$ . Thus  $X = \operatorname{Ker}(c)$  is a required object in this case.

We next come to general case. Let  $W_1$  be the maximal proper  $L$ -submodule of  $W$ . Then we see that  $(L, W/W_1 \subset V/W_1)$  satisfies the assumption of the lemma with  $W/W_1$  irreducible. By the argument above, we therefore see that there exists an  $L$ -submodule  $Y$  which contains  $W_1$  as a submodule of codimension 1 such that

$$V/W_1 = Y/W_1 \oplus W/W_1$$

holds. Since  $(L, W_1 \subset Y)$  also satisfies the assumption of the lemma with  $\dim(Y) < \dim(V)$ , we deduce by induction that the lemma holds in general. □

LEMMA 5.52. *Let  $L$  be a Lie algebra over a commutative ring  $k$ . Then for any  $L$ -modules  $V, W$ , each of the vector spaces*

$$\operatorname{Hom}_{k\text{-linear}}(V, W)$$

and

$$V \otimes_k W$$

admits a structure of  $L$ -module. Namely,

$$(x.f)(v) = x.(f(v)) - f(x.v) \quad (\forall x \in L, \forall f \in \operatorname{Hom}_{k\text{-linear}}(V, W) \forall v \in V),$$

$$(x.v \otimes_k w) = (x.v) \otimes w + v \otimes_k (x.w) \quad (\forall x \in L, \forall v \in V, \forall w \in W).$$

PROOF. Easy. □

**THEOREM 5.53** (Weyl). *Let  $k$  be a field of characteristic  $p$  (which may be 0). Let  $L$  be a non degenerate Lie algebra over  $k$ . Let  $V$  be a finite dimensional  $L$ -module. Assume:*

- (1)  $p \in \text{Ccs}(\dim(V)^2)$ .
- (2)  $p > \dim(V)^2$  or  $p = 0$ .

*Then for any  $L$ -submodule  $W$  of  $V$ , The exact sequence*

$$0 \rightarrow W \rightarrow V \rightarrow V/W \rightarrow 0$$

*of  $L$ -modules splits.*

**PROOF.** Let  $W$  be an  $L$ -submodule of  $V$ . Let us define the following  $L$ -modules.

$$V_1 = \{f \in \text{Hom}_{k\text{-linear}}(V, W); f|_W \in k \cdot 1_W\}$$

$$W_1 = \{f \in \text{Hom}_{k\text{-linear}}(V, W); f|_W = 0\}$$

Then it is easy to see that the triple  $(L, W_1 \subset V_1)$  satisfies the assumption of Lemma 5.51. We therefore have an element  $f \in \text{Hom}_{k\text{-linear}}(V, W)$  which satisfies the following conditions.

- (1)  $f|_W \in k \cdot 1_W$ .
- (2)  $f|_W \neq 0$ .
- (3)  $x.f = 0 \quad (\forall x \in L)$  (In other words,  $f$  is a  $L$ -linear homomorphism).

by dividing by a suitable element in  $k$ , we may assume  $f|_W = 1_W$ . Then  $f$  gives a splitting of the exact sequence

$$0 \rightarrow W \rightarrow V \rightarrow V/W \rightarrow 0.$$

□

### 5.13. Semi direct products of Lie algebras.

**DEFINITION 5.54.** Let  $L$  be a Lie algebra over a commutative ring  $k$ . Then:

- (1) A ( $k$ -linear) **derivation** of  $L$  is a  $k$ -linear map  $D : L \rightarrow L$  such that it obeys the following ‘‘Leibniz rule’’.

$$D([x, y]) = [Dx, y] + [x, Dy] \quad (\forall x, y \in L).$$

- (2) We denote by  $\text{Der}_k(L)$  the set of all derivations of  $L$ .

**LEMMA 5.55.** (1) *Any derivation  $D$  of a Lie algebra  $L$  is lifted to a derivation on the universal enveloping algebra  $U(L)$ .*

- (2)  $\text{Der}_k(L)$  forms a Lie algebra under the usual  $k$ -linear structure and the usual commutator as the bracket product.

DEFINITION 5.56. Let  $L_1, L_2$  be Lie algebras over a commutative ring  $k$ . we say “ $L_1$  **acts on**  $L_2$  **as a derivation**” if there is given a Lie algebra homomorphism

$$\pi : L_1 \rightarrow \text{Der}_k(L_2).$$

If the action  $\pi$  is obvious in context, we shall simply denote  $x.y$  instead of  $\pi(x).y$ .

DEFINITION 5.57. Let  $L_1, L_2$  be Lie algebras over a commutative ring  $k$ . Assume there is given an action  $\pi$  of  $L_1$  on  $L_2$ . Then we define a **semi direct product**  $L_1 \ltimes_{\pi} L_2$  of  $L_1$  and  $L_2$  by introducing the  $k$ -module  $L_1 \oplus L_2$  with the following bracket product.

$$[(x_1, x_2), (y_1, y_2)] = ([x_1, y_1], [x_2, y_2] + x_1.y_2 - y_1.x_2).$$

Note that

- (1)  $L_1$  and  $L_2$  are (identified with) subalgebras of  $L_1 \ltimes_{\pi} L_2$ .
- (2) Further more,  $L_2$  is an ideal of  $L_1 \ltimes_{\pi} L_2$ .
- (3) For  $x \in L_1$  and  $y \in L_2$ , we have

$$[x, y]_{L_1 \ltimes L_2} = x.y.$$

#### 5.14. Levi decomposition.

DEFINITION 5.58. Let  $L$  be a Lie algebra over a field  $k$ . Let  $R$  be the radical of  $L$ . A **Levi-subalgebra** of  $L$  is a subalgebra  $L_1$  of  $L$  such that  $L$  is a direct sum of  $L_1$  and  $R$  as a vector space over  $k$ .

We have the following obvious lemma.

LEMMA 5.59. *Let  $L$  be a Lie algebra over a field  $k$ . Let  $R$  be the radical of  $L$ . Let  $L_1$  be a Levi-subalgebra of  $L$ . Then:*

- (1)  $L \cong L_1 \ltimes R$ .
- (2)  $L_1 \cong L/R$ .

*In particular the isomorphism class of  $L_1$  is unique.*

□

LEMMA 5.60. *Let  $n$  be a positive integer. Let  $L$  be an  $n$ -dimensional Lie algebra over a field  $k$  of characteristic  $p \in \text{Ccs}(n^2)$ . Assume*

- (1) *The radical  $R$  of  $L$  is abelian.*
- (2) *The center of  $L$  is trivial.*
- (3)  *$p > n^2$  or  $p = 0$ .*

*Then:*



(1)

$$V = \text{Hom}_{k\text{-linear}}(L, R)$$

admits an action  $\alpha$  of  $L$ . Namely,

$$(\alpha(x).\varphi)(y) = [x, \varphi(y)] - \varphi([x, y]) \quad (x \in L, y \in L).$$

(2)

$$V_1 = \{\varphi \in \text{Hom}_{k\text{-linear}}(L, R); \varphi|_R \in k \cdot \text{id}_R\}$$

is an  $L$ -submodule of  $V$ .

(3)

$$V_0 = \{\varphi \in \text{Hom}_{k\text{-linear}}(L, R); \varphi|_R = 0\}$$

is an  $L$ -submodule of  $V_1$ .

(4)

$$U = \text{ad}_L(R)$$

is an  $L$ -submodule of  $V_0$ .

(5)  $R.V_1 \subset U$ .(6) There exists  $\chi \in \text{Hom}_{k\text{-linear}}(L, R)$  such that

$$\chi|_R = \text{id}_R, \quad L.\chi \subset U$$

holds.

(7) For any  $x \in L$ , there exists a unique element  $r_x \in R$  such that

$$x.\chi = \text{ad } r_x$$

holds.

(8)

$$h : L \ni x \mapsto -r_x \in R$$

is a  $k$ -linear projection.

(9)  $L_1 = \text{Ker}(h)$  is a Levi subalgebra of  $L$ .

PROOF. (1):follows from the general theory.

(2),(3),(4):follows easily from the definition of  $\alpha$ .

(5): For any  $r \in R$  and for any  $\varphi \in V_1$ , we have

$$r.\varphi = (z \mapsto [r, [\varphi(z)] - \varphi([r, z]]) = (z \mapsto 0 - c_\varphi[r, z]) = \text{ad}_L(-c_\varphi r) \in U.$$

(6):  $V_1/U$  is an  $L/R$ -module which has  $V_0/U$  as a submodule of codimension 1. Thus by using Lemma 5.51 (Weyl's theorem on irreducibility (codimension 1 case)), we see that there exists a 1-dimensional  $L/R$ -submodule  $X/U$  of  $V_1/U$  (where  $X$  is a  $L$  submodule of  $L$ ) such that

$$V_1/U = V_0/U \oplus X/U$$

holds. Since  $X/U$  is 1 dimensional, the action of the semisimple Lie algebra  $L/R$  on  $X/U$  is trivial. Thus we see that there exists an element  $\chi \in \text{Hom}_{k\text{-linear}}(L, R)$  such that

$$\chi|_R = \text{id}_R, \quad L \cdot \chi \subset U$$

holds.

(7): Since  $x \cdot \chi$  belongs to  $R$ , we know the existence of  $r_x$ . The uniqueness of the  $r_x$  follows from the assumption that the center of  $L$  is trivial.  
 (8),(9): easy. □

LEMMA 5.61. *Let  $n$  be a positive integer. Let  $L$  be an  $n$ -dimensional Lie algebra over a field  $k$  of characteristic  $p \in \text{Ccs}(n^2)$ . Assume*

- (1) *The radical  $R$  of  $L$  is abelian.*
- (2)  *$p > n^2$  or  $p = 0$ .*

*Then  $L$  has a Levi subalgebra.*

PROOF. Let  $Z$  be the center of  $L$ . Applying the previous lemma to  $L/Z$ , We see that there exists an Levi subalgebra  $L_0/Z$  of  $L/Z$ . Now

$$0 \rightarrow Z \rightarrow L_0 \rightarrow L_0/Z \rightarrow 0$$

is a short exact sequence of  $L_0/Z$ -module, and so it therefore splits. (Theorem 5.53 (Weyl's theorem of irreducibility.)) Thus  $L_0$  has a subalgebra  $L_1$  which is stable under action of  $L_0/Z$ . That means,  $L_1$  is a Levi subalgebra of  $L_0$ . So  $L_1$  is also a Levi subalgebra of  $L$ . □

THEOREM 5.62 (Levi decomposition of a Lie algebra). *Let  $n$  be a positive integer. Let  $p \in \text{Ccs}(n^2)$ . Let  $L$  be a  $n$ -dimensional Lie algebra over a field of characteristic  $p$ . Then  $L$  has a Levi subalgebra  $L_1$ . In other words,  $L$  may be expressed as a semi direct product*

$$L = L_1 \ltimes R$$

*where  $L_1$  is a semisimple (Levi) subalgebra of  $L$ , and  $R$  is a solvable (radical) ideal of  $L$ .*

PROOF. If  $R = 0$ , then we only need to set  $L_1 = L$ . So let us assume  $R \neq 0$ . Let us put

$$R_1 = [R, R].$$

Then from the definition, we  $R/R_1$  is an abelian Lie algebra. It is also easy to verify that  $R_1$  is an ideal of  $L$ . ( $R_1$  is a **characteristic ideal** of  $R$ .) We apply the preceding lemma for  $R/R_1 \subset L/R_1$  to obtain a Levi subalgebra  $M/R_1$  of  $L/R_1$ . Then  $M$  satisfies the following relations.

$$M + R = L, \quad M \cap R = R_1.$$

Since  $R$  is solvable (and we have assumed  $R \neq 0$ ), we see that  $\dim(M)$  is strictly smaller than  $\dim(L)$ . By induction  $M$  have a Levi subalgebra  $M_1$ . Then it is clear that  $M_1$  is a Levi subalgebra of  $L$ . □

**5.15. Abstract Jordan Chevalley decomposition.**

**PROPOSITION 5.63.** *Let  $n$  be a positive integer. Let  $L$  be a  $n$ -dimensional semisimple Lie algebra over a field  $k$  of characteristic  $p \in \text{Ccs}(n^4)$ . Then any derivation  $D \in \text{Der}_k(L)$  of  $L$  is inner. That is, there exists an element  $x = x_D$  such that*

$$D(y) = [x_D, y] = \text{ad}(x_D)(y).$$

**PROOF.**  $\text{Der}_k(L)$  is itself a Lie algebra. Sending each element  $x$  of  $L$  to its “inner derivation”  $\text{ad}(x)$ , we obtain a Lie algebra homomorphism

$$\text{ad} : L \rightarrow \text{Der}_k(L)$$

We note that  $\dim(\text{Der}_k(L)) \leq \dim(L)^2$ , and that  $\text{ad}$  may be viewed as a homomorphism of  $L$ -modules. ( $L$  acts on  $\text{Der}_k(L)$  via  $\text{ad}$ . Namely,

$$x.D = \text{ad}(x).D = [\text{ad}(x), D] = [x, D\bullet] - D([x, \bullet]) = -\text{ad}(D.x)$$

holds for any  $x \in L$  and for any  $D \in \text{Der}_k(L)$ .) By the Weyl’s theorem on complete reducibility, we see that there exists a direct sum decomposition

$$\text{Der}_k(L) = \text{ad}(L) \oplus X$$

of  $L$ -modules. Then for any  $D \in X$  and for any  $x \in L$ , we see that

$$x.D (= -\text{ad}(D.x)) \in X \cap \text{ad}(L) = 0.$$

So  $D = 0$ . That means,  $X = 0$ . □

**PROPOSITION 5.64.** *Let  $n$  be a positive number Let  $k$  be a separably closed field of characteristic  $p \in \text{Ccs}(n^4)$ . We assume further that  $n$  is invertible in  $k$ . (This assumption is provided just in case: it probably is not necessary because the assumption  $p \in \text{Ccs}(n^4)$  is presumably much stronger.) Let  $L \subset \mathfrak{gl}_n(k)$  be a linear semisimple Lie algebra. We assume that the representation  $L \curvearrowright k^n$  is irreducible. Then for any element  $x \in L$ , its semisimple part  $x_s$  and its nilpotent part  $x_n$  in  $\mathfrak{gl}_n(k)$  lies in  $L$ .*

**PROOF.** We may assume  $k$  is algebraically closed. Let  $x \in L$  It is enough to prove  $x_n \in L$ . There exists a polynomial  $f \in k[X]$  such that  $x_n = f(x)$ . Thus we see

$$\text{ad } x_n(L) \subset L.$$

Thus  $\text{ad } x_n$  is a derivation of  $L$ . By the preceding lemma we see that there exists an element  $y \in L$  such that

$$\text{ad } x_n = \text{ad } y$$

By Schur's lemma, we see that there exists a constant  $c \in k$  such that

$$x_n = y + c \cdot 1_n.$$

Let us compute traces of both hand sides. Since  $L = [L, L]$  ( $L$  has no non-trivial ideals.), we have  $\text{tr}(y) = 0$ . Since  $x_n$  is nilpotent, we have  $\text{tr}(x_n) = 0$ . Thus we conclude  $c = 0$  (as we assumed  $n$  is invertible in  $k$ .)  $\square$

**PROPOSITION 5.65.** *Let  $n$  be a positive integer. Let  $L$  be a semisimple Lie algebra over a separably closed field  $k$  of characteristic  $p \in \text{Ccs}(n^4)$ . Let  $V_1 = (V_1, \pi_1), V_2 = (V_2, \pi_2)$  be faithful irreducible representations of  $L$  with dimensions less than or equal to  $n$ . Then for any  $x \in L$ , the Jordan Chevalley decomposition of  $x$*

$$x = x_s^{(1)} + x_n^{(1)}$$

*with respect to  $V_1$  and that*

$$x = x_s^{(2)} + x_n^{(2)}$$

*with respect to  $V_2$  coincides.*

**PROOF.** We consider a faithful representation  $(V, \pi) = (V_1 \oplus V_2, \pi_1 \oplus \pi_2)$ . For any  $x \in L$ ,

$$\pi(x) = \begin{pmatrix} \pi_1(x) & 0 \\ 0 & \pi_2(x) \end{pmatrix} = \begin{pmatrix} \pi_1(x_s^{(1)}) & 0 \\ 0 & \pi_2(x_s^{(2)}) \end{pmatrix} + \begin{pmatrix} \pi_1(x_n^{(1)}) & 0 \\ 0 & \pi_2(x_n^{(2)}) \end{pmatrix}$$

satisfies the requirement for the Jordan Chevalley decomposition so by the uniqueness we see

$$\pi(x)_s = \begin{pmatrix} \pi_1(x_s^{(1)}) & 0 \\ 0 & \pi_2(x_s^{(2)}) \end{pmatrix}, \quad \pi(x)_n = \begin{pmatrix} \pi_1(x_n^{(1)}) & 0 \\ 0 & \pi_2(x_n^{(2)}) \end{pmatrix}.$$

Now we argue in a same way as in the proof of the previous proposition and see that there exists a unique element  $y \in L$  such that

$$\text{ad}(x_n) = \text{ad}(y)$$

holds. By comparing entries, we obtain

$$\text{ad}(\pi_1(y)) = \text{ad}(\pi_1(x_n^{(1)})), \quad \text{ad}(\pi_2(y)) = \text{ad}(\pi_2(x_n^{(2)})).$$

Since  $L$  has trivial center, we have

$$\pi_1(y) = \pi_1(x_n^{(1)}), \quad \pi_2(y) = \pi_2(x_n^{(2)}).$$

Thus  $y = x_n^{(1)} = x_n^{(2)}$ . □

**DEFINITION 5.66.** Let  $n$  be a positive integer. Let  $L$  be an  $n$ -dimensional semisimple Lie algebra over a separably closed field  $k$  of characteristic  $p \in \text{Ccs}(n^4)$ . Then the **abstract Jordan Chevalley decomposition** of  $x$  is an decomposition

$$x = x_s + x_n \quad (x_s, x_n \in L)$$

such that

$$\text{ad}(x) = \text{ad}(x_s) + \text{ad}(x_n)$$

is the Jordan Chevalley decomposition.

**PROPOSITION 5.67.** *Let  $n$  be a positive integer. Let  $L$  be an  $n$ -dimensional Lie algebra over a separably closed field  $k$  of characteristic  $p \in \text{Ccs}(n^4)$ . Then the abstract Jordan Chevalley decomposition of  $x$  exists. If furthermore there is given a  $m$ -dimensional representation  $(V, \pi)$  of  $L$  and  $p \in \text{Ccs}(m^4)$ , then*

$$\pi(x) = \pi(x_s) + \pi(x_n)$$

*gives the Jordan Chevalley decomposition of  $x$ .*

**PROOF.** Easy exercise. (Be sure to use Weyl's theorem of complete reducibility. By taking quotient by a certain ideals (kernels of representations) one may reduce the proposition to a case where  $L$  is semisimple and  $\pi$  is faithful and irreducible. ) □

## REFERENCES

- [1] J. E. Humphreys, *Introduction to Lie algebras and representation theory*, Springer-Verlag, 1972.