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LINEAR ALGEBRA 2

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MODULE SUMMARY. This lecture notes are based on a course at the University of Leeds.

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1. VECTOR SPACES AND SUBSPACES

Vector spaces have two built-in concepts.

- (1) Vectors: these can be added, subtracted from each other.
- (2) Scalars: these can be added, subtracted, multiplied, divided (except by zero).

Also, we can multiply a vector by a scalar (and get a vector). The set of scalars in this course is usually \mathbb{R} (the real numbers) or \mathbb{C} (the complex numbers), but in general there are many other possibilities.

Example 1.1. Here are some vector spaces. The term 'vector space' really refers to the set of vectors.

(1) \mathbb{R}^3 (or more generally, \mathbb{R}^n , for any fixed $n \ge 1$). Here the set of scalars is \mathbb{R} . Observe

$$\begin{array}{rcl} (x_1,x_2,x_3)+(y_1,y_2,y_3)&=&(x_1+y_1,x_2+y_2,x_3+y_3)\\ (x_1,x_2,x_3)-(y_1,y_2,y_3)&=&(x_1-y_1,x_2-y_2,x_3-y_3),\\ &a(x_1,x_2,x_3)&=&(ax_1,ax_2,ax_3), \end{array}$$

for any $x_i, y_i, a \in \mathbb{R}$.

- (2) \mathbb{C}^n . This is similar, but here the scalars are complex numbers.
- (3) M_{m,n}(ℝ), the set of all m × n (i.e. m rows, n columns) matrices over ℝ (the set of scalars). Here the vectors ARE matrices, and vector addition is matrix addition. M_{m,n}(ℂ), the set of all m × n-matrices over ℂ (and the set of scalars is ℂ).
- (4) $P_n(t)$, the set of all polynomials in variable t of degree at most n, i.e.

$$\mathbf{P}_{\mathbf{n}}(\mathbf{t}) \coloneqq \{\mathbf{a}_0 + \mathbf{a}_1 \mathbf{t} + \ldots + \mathbf{a}_n \mathbf{t}^n : \mathbf{a}_i \in \mathbb{R}\}.$$

The scalars are \mathbb{R} .

(5) Fix a set X. The set \mathbb{R}^X is the collection of all functions from X into \mathbb{R} . For f, $g \in \mathbb{R}^X$ and $x \in X$, we have

$$(f+g)(x) = f(x) + g(x)$$

$$(af)(x) = af(x).$$

Likewise, \mathbb{C}^X (functions from X to \mathbb{C}) is a vector space with scalars \mathbb{C} .

1.1. Axioms for a vector space. I'll generally use u, v, w, \dots etc. for vectors, and a, b, c, \dots for scalars. Let K denote the set of scalars, which in this module is always \mathbb{R} or \mathbb{C} . (In general, K can be any *field*, but I am not defining this; K can even be finite.) Let V be a non-empty set of *vectors*, and suppose there are rules of addition of vectors, and scalar multiplication, so that if $u, v \in V$ then the sum $u + v \in V$, and if $a \in K, v \in V$, then $av \in V$. Then V is called a *vector space* over K if the following axioms hold.

- (1) V is an abelian group under +, i.e.
 - (a) $\forall u, v \in V \ (u + v \in V)$ (V is closed under vector addition).
 - (b) (u + v = v + u), $\forall u, v \in V$ (+ is *commutative*, i.e. V is an *abelian* group with operation +).
 - (c) $(\mathbf{u} + \mathbf{v}) + \mathbf{w} = \mathbf{u} + (\mathbf{v} + \mathbf{w}) \forall \mathbf{u}, \mathbf{v}, \mathbf{w} \in \mathbf{V}$, (+ is associative);
 - (d) there is a zero vector $0 \in V$ such that $\forall v \in V (0 + v = v + 0 = v)$;
 - (e) for all $\nu \in V$ there is an 'inverse vector' $-\nu$ in V such that $\nu + (-\nu) = (-\nu) + \nu = 0$;
- (2) Scalar multilcation satisfies (for all $u, v \in V$ and $a, b \in K$):
 - (a) $\forall u \in V, \forall a \in \mathbb{K} \ (av \in V)$ (V is closed under scalar multiplication).
 - (b) a(u + v) = au + av;
 - (c) (a + b)u = au + bu;
 - (d) (ab)u = a(bu);
 - (e) 1u = u.

In practice, we don't usually have to check all these, to see that something is a vector space.

Remark 1.2. Note that sometimes 0 denotes the zero vector, sometimes the zero scalar – this should be clear from context.

The next theorem collects facts really obvious for \mathbb{R}^n or \mathbb{C}^n . We need to know them for *all* vector spaces, i.e. that they follow from the axioms for vector spaces.

Theorem 1.3. Let V be a vector space over a field K. Then

- (1) *if* $a \in K$ and $0 \in V$ then $a \cdot 0 = 0$;
- (2) *if* $0 \in K$ and $v \in V$ then $0 \cdot v = 0$;
- (3) if av = 0 then either a = 0 or v = 0;
- (4) (-a)v = a(-v) = -(av).

Proof. I'll just do (i), (iii).

- (i) $\mathbf{a} \cdot \mathbf{0} = \mathbf{a} \cdot (\mathbf{0} + \mathbf{0}) = \mathbf{a} \cdot \mathbf{0} + \mathbf{a} \cdot \mathbf{0}$. Now add $-(\mathbf{a} \cdot \mathbf{0})$ to both sides. We find $\mathbf{0} = \mathbf{a} \cdot \mathbf{0} + \mathbf{0} = \mathbf{0} + \mathbf{a} \cdot \mathbf{0}$.
- (iii) Suppose that av = 0 and $a \neq 0$. Then there is $a^{-1} \in K$, so $a^{-1}av = a^{-1} \cdot 0 = 0$. But also, $a^{-1}av = 1 \cdot v = v$. Hence v = 0.

Example 1.4. To demonstrate that choice of axioms is a delicate matter we will consider examples which violet these "obvious" properties.

On the set \mathbb{R}^2 of ordered pairs (x, y) define + as usual, i.e. (x, y) + (u, v) = (x + u, y + v), but define the multiplication by scalar as $\lambda \cdot (x, y) = (x, 0)$.

Since the addition is defined in the usual way we could easily check all corresponding axioms:

- (1) $a + b \in V$.
- (2) a + b = b + a.
- (3) a + (b + c) = (a + b) + c.
- (4) $\mathbf{a} + \mathbf{0} = \mathbf{a}$, where $\mathbf{0} = (0, 0)$.

(5) For any vector $\mathbf{a} = (x, y)$ there is $-\mathbf{a} = (-x, -y)$ such that $\mathbf{a} + (-\mathbf{a}) = \mathbf{0}$.

However for our strange multiplication $\lambda \cdot (x, y) = (x, 0)$ we should be more careful:

- (1) $\lambda \cdot \mathbf{a} \in V$.
- (2) $\lambda \cdot (\mathbf{a} + \mathbf{b}) = \lambda \cdot \mathbf{a} + \lambda \cdot \mathbf{b}$, because $\lambda \cdot ((x, y) + (u, v)) = \lambda \cdot (x + u, y + v) = (x + u, 0)$ $\lambda \cdot (x, y) + \lambda \cdot (u, v) = (x, 0) + (u, 0) = (x + u, 0).$
- (3) $(\lambda + \mu) \cdot \mathbf{a} \neq \lambda \cdot \mathbf{a} + \mu \cdot \mathbf{a}$, because $(\lambda + \mu) \cdot (x, y) = (x, 0)$ however $\lambda \cdot (x, y) + \mu \cdot (x, y) = (x, 0) + (x, 0) = (2x, 0)$.
- (4) $(\lambda \mu) \cdot \mathbf{a} = \lambda(\mu \cdot \mathbf{a}) = \mu(\lambda \cdot \mathbf{a})$, i.e. the *associative law* holds (why?)
- (5) $1 \cdot \mathbf{a} \neq \mathbf{a}$ (why?)

This demonstrate that the failing axioms could not be derived from the rest, i.e they are independent.

1.2. **Subspaces.** Let *W* be a subset of a vector space *V* over *K*. Then *W* is a *subspace* of *V* if *W* is also a vector space over *K*. This looks boring, with lots of axioms to check, but it's easy to check in practice, because of the following.

Theorem 1.5. $W \subseteq V$ is a subspace of V if and only if

- (1) $0 \in W$,
- (2) W is closed under addition of vectors, i.e. if $v, w \in W$ then $v + w \in W$;
- (3) W is closed under scalar multiplication, i.e. if $w \in W$ and $a \in K$ then $aw \in W$.

Proof. Clearly, if W is a subspace then 1–3 hold. We show the converse, so assume 1–3 above. Then W is a group: the inverse of w is $-w = -(1 \cdot w) = (-1)w$, so lies in W.

The other properties all hold in W since they hold in V, e.g. w + 0 = 0 + w = w holds if $w \in W$, since $W \subseteq V$ so $w \in V$, and the above rule holds in V.

Corollary 1.6. $W \subseteq V$ is a subspace of V if and only if

- (1) $0 \in W$;
- (2) $av + bw \in W$ whenever $a, b \in K$ and $v, w \in W$.

Example 1.7. (1) Always, {0} and V are subspaces of V;

- (2) In the vector space \mathbb{R}^2 , the set of points on a line through the origin (eg 2x+3y = 0) is a subspace, but 2x + 3y = 1 is not.
- (3) In $M_{n,n}(\mathbb{R})$, consider the set X of *symmetric* $n \times n$ matrices. Here $A = (a_{ij})$ is *symmetric* if $a_{ij} = a_{ji}$ for all i, j, i.e. if $A = A^T$. An example of a symmetric matrix $\begin{pmatrix} 1 & 2 & 4 \end{pmatrix}$
 - is $\begin{pmatrix} 1 & 2 & 4 \\ 2 & 3 & 6 \\ 4 & 6 & 5 \end{pmatrix}$. Then X is a subspace of $M_{n,n}(\mathbb{R})$.
- (4) The solution space of any system of homogeneous linear equations in n variables (coefficients in ℝ) is a subspace of ℝⁿ.
- *Remark* 1.8. (1) If U_1, U_2 are subspaces of V then $U_1 \cap U_2$ is a subspace of V. For $0 \in U_1 \cap U_2$, and if $u, v \in U_1 \cap U_2$, and $a, b \in K$, then $au + bv \in U_1$, $au + bv \in U_2$, so $au + bv \in U_1 \cap U_2$, so Corollary 1.6 applies.
 - (2) In general, the *union* of two subspaces is not a subspace; for if $u_1 \in U_1$ and $u_2 \in U_2$, then possibly $u_1 + u_2 \notin U_1 \cup U_2$. For example, in \mathbb{R}^2 , let

$$\begin{array}{rcl} U_1 &:= & \{(x,0): x \in \mathbb{R}\}, \\ U_2 &:= & \{(0,y): y \in \mathbb{R}\}. \end{array}$$

Then $(1,0), (0,1) \in U_1 \cup U_2$, but $(1,1) \notin U_1 \cup U_2$.

Definition 1.9. Given a set S of vectors in a vector space V, a vector v is a *linear combination* of elements of S if $v = a_1s_1 + \ldots + a_ks_k$ for some $s_1, \ldots, s_k \in S$ and $a_1, \ldots, a_k \in K$. We denote by span(S) the set of linear combinations of elements of S.

Theorem 1.10. Let S be a non-empty subset of V. Then

- (1) $\operatorname{span}(S)$ is a subspace of V;
- (2) any subspace of V which contains S also contains span(S).

Proof. As usual, we apply Corollary 1.6.

(1) Clearly $0 \in \text{span}(S)$, e.g. as $0 = 0 \cdot s_1$ for any $s_1 \in S$. Let $v, w \in \text{span}(S)$, say $v = a_1s_1 + \ldots + a_ks_k$, $w = b_1t_1 + \ldots + b_kt_k$ (so $a_i, b_i \in K$ and $s_i, t_i \in S$). Suppose $a, b \in K$. Then av + bw =

 $a(a_1s_1 + \ldots + a_ks_k) + b(b_1t_1 + \ldots + b_kt_k) = (aa_1)s_1 + \ldots + (aa_k)s_k + (bb_1)t_1 + \ldots + (bb_k)t_k.$

(2) This is clear, (formally, by induction).

Example 1.11. In \mathbb{R}^2 , the smallest subspace containing (1, 1), (2, 3) is all of \mathbb{R}^2 , as any vector (x, y) can be written as $\mathfrak{a}(1, 1) + \mathfrak{b}(2, 3)$, as we can always solve

$$\begin{array}{rcl} a+2b & = & x, \\ a+3b & = & y. \end{array}$$

Thus, span{(1, 1), (2, 3)} = \mathbb{R}^2 .

1.3. Direct Sums.

Definition 1.12. The *sum* of two subspaces U_1 , U_2 of V is the set

 $U_1 + U_2 := \{u_1 + u_2 : u_1 \in U_1, u_2 \in U_2\}.$

Theorem 1.13. *If* U_1 , U_2 *are subspaces of* V*, then* $U_1 + U_2$ *is a subspace of* V*, and is the smallest subspace containing* U_1 , U_2 .

Proof. (1) $U_1 + U_2$ is a subspace. Indeed, $0 = 0 + 0 \in U_1 + U_2$, and if $u_1, u'_1 \in U_1$, $u_2, u'_2 \in U_2$, then

 $a(u_1+u_2)+b(u_1'+u_2')=(au_1+bu_1')+(au_2+bu_2')\in U_1+U_2.$

(2) Any subspace of V which contains U_1 , U_2 must contain all the vectors $u_1 + u_2$, so contains $U_1 + U_2$.

Remark 1.14. We have $U_1 + U_2 = \text{span}(U_1 \cup U_2)$.

Example 1.15. Work in \mathbb{R}^3 . Let

$$\begin{array}{rcl} V_0 & := & \{(x,y,z): x+y+z=0\}, \\ V_1 & := & \{a(1,0,0): a\in \mathbb{R}\}, \\ V_2 & := & \{a(1,0,-1): a\in \mathbb{R}\}. \end{array}$$

Then $V_0 + V_1 = \mathbb{R}^3$, since for $a, b, c \in \mathbb{R}$,

 $(a,b,c) = (-b-c,b,c) + (a+b+c,0,0) \in V_0 + V_1.$

Also $V_0 + V_2 = V_0$, since $V_2 \subset V_0$ (so anything in $V_0 + V_2$ already lies in V_0). The subspace $V_1 + V_2$ is a plane.

Definition 1.16. If U_1 , U_2 are subspaces of V, and $U_1 \cap U_2 = \{0\}$, then $U_1 + U_2$ is called a *direct sum* of U_1 , U_2 , written $U_1 \oplus U_2$.

Remark 1.17. Any vector in $U_1 \oplus U_2$ has the form $u_1 + u_2$ (where $u_i \in U_i$) in a *unique* way. For suppose $u_1 + u_2 = u'_1 + u'_2$ (where $u_i, u'_i \in U_i$). Then $u_1 - u'_1 = u_2 - u'_2$. But the left hand side lies in U_1 , and the right hand side lies in U_2 , so $u_1 - u'_1 \in U_1 \cap U_2 = \{0\}$, so $u_1 - u'_1 = 0$, so $u_1 = u'_1$, and likewise $u_2 = u'_2$.

Example 1.18. (1) In Example 1.15, $\mathbb{R}^3 = V_0 \oplus V_1$. For if $(a, 0, 0) \in V_0$, then a+0+0=0, so a = 0.

(2) Let $V := P_2(t)$ be the set of all polynomials over \mathbb{R} of degree at most 2 (see Example 1.1). Let

$$\begin{array}{lll} V_0 &:= & \{at^2: a \in \mathbb{R}\}, \\ V_1 &:= & \{bt+c: b, c \in \mathbb{R}\}, \\ V_2 &:= & \{a(t^2+1): a \in \mathbb{R}\}. \end{array}$$

Then V_0 , V_1 , V_2 are subspaces of V (check this!). And $V = V_0 \oplus V_1 = V_2 \oplus V_1$. To see that $V = V_2 \oplus V_1$, observe that $at^2 + bt + c = a(t^2 + 1) + (bt + (c - a))$; it is trivial that $V_1 \cap V_2 = \{0\}$.

Definition 1.19. If $U_1 \oplus U_2 = V$, then U_2 is a *complement* of U_1 in V.

Remark 1.20. The complement is not necessarily unique; e.g., in Example 1.18, V_0 and V_2 are both complements of V_1 .

2. BASES AND DIMENSION

The material in this section (up to and including 2.8) was done for subspaces of \mathbb{R}^n in MATH1011.

Definition 2.1. A set S *spans* the vector space V if span(S) = V, i.e., any $v \in V$ can be written $v = a_1s_1 + \ldots + a_ns_n$ for some $a_i \in K$, $s_i \in S$.

The set S is *linearly independent* if the only linear combination of members of S giving 0 is the trivial one, i.e., if $0 = a_1s_1 + ... + a_ns_n$ (where $s_i \in S$ are distinct), then $a_1 = ... = a_n = 0$.

Example 2.2. Some examples in $V = \mathbb{R}^2$.

- (1) (1, 1), (2, 3) span \mathbb{R}^2 and are linearly independent.
- (2) (1,1), (2,2) span \mathbb{R}^2 but are not linearly independent.
- (3) (1,1), (2,2), (2,3) span \mathbb{R}^2 but are not linearly independent.
- (4) (1,1) is linearly independent but does not span \mathbb{R}^2 .

Exercise 2.3. (1) Any set including the zero vector is linearly independent.

- (2) Any subset of a linearly independent subset is linearly independent as well.
- (3) Any set, which includes a spanning set, is spanning as well.
- (4) A spanning set expanded by a single vector become linearly dependent.
- (5) A linearly independent set with any single vector removed cannot be spanning.

To check whether a set of r vectors in \mathbb{R}^n is linearly independent, form the $r \times n$ matrix whose rows are these vectors, and put it into row reduced form. The original set is linearly independent if and only if the row reduced form has r non-zero rows.

Definition 2.4. A *basis* of V is a set of vectors in V which is both linearly independent and spans V.

Remark 2.5. The set $\{s_1, \ldots, s_n\}$ is a basis of V if and only if every vector $v \in V$ can be written as $v = a_1s_1 + \ldots + a_ns_n$ ($a_i \in K$, s_i distinct members of S) in a *unique* way (up to order). To see this uniqueness, note that if $\{s_1, \ldots, s_n\}$ is linearly independent and $a_1s_1 + \ldots + a_ns_n = b_1s_1 + \ldots + b_ns_n$, then $0 = (a_1 - b_1)s_1 + \ldots + (a_n - b_n)s_n$, so $a_i = b_i = 0$ for all i, so the expressions are the same.

If $\{s_1, \ldots, s_n\}$ is a basis for V then V 'looks like' K^n ; for each $\nu \in V$ can be written uniquely as $\nu = a_1s_1 + \ldots + a_ns_n$, so we can identify ν with $(a_1, \ldots, a_n) \in K^n$.

Theorem 2.6. Suppose that V is a vector space over K, and $S \subseteq V$ is linearly independent. Then there is a basis B of V with $S \subseteq B$.

Proof. We'll do the special case when S is finite, and there are $v_1, \ldots, v_m \in V$ with $V = \text{span}(v_1, \ldots, v_m)$. Suppose $S = \{s_1, \ldots, s_n\}$.

Case 1 span(S) = V. Now S is a basis of V, so put B := S.

Case 2 There is some $v_i \in \{v_1, \ldots, v_m\}$ not in span(S). Now $S \cup \{v_i\}$ is linearly independent. For if $a_1s_1 + \ldots + a_ns_n + bv_i = 0$, then b = 0, for otherwise $v_i = -b^{-1}(a_1s_1 + \ldots + a_ns_n) \in \text{span}(S)$. Hence, as S is linearly independent, also $a_1 = \ldots = a_n = 0$. Now repeat the above argument with $S \cup \{v_i\}$ in place of S. If we do not stop before then, we eventually obtain a linearly independent set $S \cup \{v_1, \ldots, v_m\}$, which certainly spans V so is a basis. If we stop before then, we have a basis.

Corollary 2.7. Suppose V is a vector space over K, $S \subseteq V$ is linearly independent, and $A \subseteq V$ satisfies span(A) = V. Then there is a basis B of V with $S \subseteq B \subseteq S \cup A$.

Proof. Our proof of Theorem 2.6 gave this, under the assumption that S, A are finite (put $A = \{v_1, \ldots, v_m\}$ in the theorem). We omit the proof in general.

Note that the proof of Theorem 2.6 give you a *method* for answering questions like Problem Sheet 2 Q1(a).

Corollary 2.8. Suppose S is a vector space over K and $A \subseteq V$ spans V. Then A contains a basis of V.

Proof. Apply Corollary 2.7 with $S = \emptyset$.

The next lemma is the key to showing any two bases of a vector space have the same size (so we can define *dimension*).

Lemma 2.9 (Exchange Lemma). (1) Suppose u_1, \ldots, u_n, v are vectors in a vector space V, and $v \in \text{span}(u_1, \ldots, u_{n-1}, u_n)$, but $v \notin \text{span}(u_1, \ldots, u_{n-1})$. Then $u_n \in \text{span}(u_1, \ldots, u_n)$ (2) Under the same assumptions, if $u_1, \ldots, u_{n-1}, u_n$ are linearly independent, so are u_1, \ldots, u_{n-1} .

Proof. (1) Since $v \in \text{span}(u_1, \ldots, u_n)$, there are scalars b_i such that

(1)

 $\nu = b_1 u_1 + \ldots + b_{n-1} u_{n-1} + b_n u_n.$

Now $b_n \neq 0$, as otherwise $v \in span(u_1, \dots, u_{n-1})$. So

$$u_n = b_n^{-1}(v - b_1u_1 - \ldots - b_{n-1}u_{n-1}) = b_n^{-1}v - b_n^{-1}b_1u_1 - \ldots - b_n^{-1}b_{n-1}u_{n-1}$$

so $u_n \in \operatorname{span}(u_1, \ldots, u_{n-1}, v)$.

(2) Suppose that u_1, \ldots, u_n are linearly independent, and that

 $a_1u_1 + \ldots + a_{n-1}u_{n-1} + a_nv = 0.$

Substituting from (1),

$$a_1u_1 + \ldots + a_{n-1}u_{n-1} + a_n(b_1u_1 + \ldots b_{n-1}u_{n-1} + b_nu_n) = 0$$

so

 $(a_1 + a_n b_1)u_1 + \ldots + (a_{n-1} + a_n b_{n-1})u_{n-1} + a_n b_n u_n = 0.$

As u_1, \ldots, u_n are linearly independent, all the coefficients are zero, so $a_n b_n = 0$. As $b_n \neq 0$, $a_n = 0$, so $a_1 u_1 + \ldots a_{n-1} u_{n-1} = 0$. But u_1, \ldots, u_{n-1} are linearly independent, so this forces $a_1 = \ldots = a_{n-1} = 0$.

Theorem 2.10. Suppose that A, B are both bases of the vector space V. Then A and B have the same number of elements.

 \square

 \square

Proof. For convenience, we assume that A, B are finite, and that A is at least as big as B. Let $B = \{v_1, \ldots, v_n\}$, and choose distinct $u_1, \ldots, u_n \in A$. We shall show that $A = \{u_1, \ldots, u_n\}$.

Now $u_1 \in V = \text{span}(v_1, \dots, v_n)$, so $u_1 = a_1v_1 + \dots + a_nv_n$ for some scalars a_i . Not all the a_i are zero, as otherwise $u_1 = 0$, contradicting that A is linearly independent. So, without loss, suppose $a_1 \neq 0$. Now $u_1 \notin \text{span}(v_2, \dots, v_n)$, as otherwise we'd have $u_1 = 0 \cdot v_1 + b_2v_2 + \dots + b_nv_n = a_1v_1 + \dots + a_nv_n$, contradicting that B is linearly independent.

So by Lemma 2.9, $v_1 \in \text{span}(u_1, v_2, ..., v_n)$, and (using 2.9(ii), $\{u_1, v_2, ..., v_n\}$ is a basis of V. Next, we apply the above argument to u_2 (in place of u_1), then to u_3 , and so on. Eventually, we find that $\{u_1, ..., u_n\}$ is a basis of V, and so $A = \{u_1, ..., u_n\}$.

Definition 2.11. The number of elements of a basis of V (which by Theorem 2.10 does not depend on the choice of basis) is called the *dimension* of V, denoted $\dim(V)$.

Corollary 2.12. Let V be a vector space and $n = \dim(V)$. If a set S has exactly n vectors then the following are equivalent:

- (1) S is linearly independent;
- (2) S spans V;
- (3) S is basis of V.

Example 2.13. (1) \mathbb{C}^n , a vector space over \mathbb{C} , has a *standard basis* $(1, 0, \dots, 0)$, $(0, 1, 0, \dots, 0)$ $(0, \dots, 0, 1)$, so has dimension n.

- (2) $P_n(t)$ has basis 1, t, t^2, \ldots, t^n , so has dimension n + 1.
- (3) $M_{m,n}(\mathbb{R})$ has a basis { $E_{ij} : 1 \le i \le m, 1 \le j \le n$ }, so has dimension mn. Here, E_{ij} is the $m \times n$ matrix with a 1 in the (i, j)-entry, zeros elsewhere.
- (4) The set of points on the line x + 2y = 0 can be described as $\{a(1, \frac{-1}{2}) : a \in \mathbb{R}\}$, so has basis $(1, \frac{-1}{2})$, so dimension 1.
- (5) The plane x + 2y + 3z = 0 in \mathbb{R}^3 has solution set

 $\{(-2a-3b, a, b) : a, b \in \mathbb{R}\} = \{a(-1, 1, 0) + b(-3, 0, 1) : a, b \in \mathbb{R}\},\$

so basis (-2, 1, 0), (-3, 0, 1), and dimension 2. Suppose we want to expand this set to a basis of \mathbb{R}^3 . Try to add in the vectors (1, 0, 0), (0, 1, 0), (0, 0, 1), and the first time the dimension goes up to 3, we have a basis. In fact, (-2, 1, 0), (-3, 0, 1), (1, 0, 0) is a basis of \mathbb{R}^3 .

Finally, an important theorem.

Theorem 2.14. Let U, V be subspaces of the finite dimensional vector space W. Then

$$\dim(\mathbf{U}) + \dim(\mathbf{V}) = \dim(\mathbf{U} \cap \mathbf{V}) + \dim(\mathbf{U} + \mathbf{V}).$$

Proof. Let w_1, \ldots, w_m be a basis for $U \cap V$. Using Theorem 2.6, we can extend this set to a basis $w_1, \ldots, w_m, u_1, \ldots, u_r$ for U, and to a basis $w_1, \ldots, w_m, v_1, \ldots, v_s$ for V. Now $\dim(U \cap V) = m, \dim(U) = m + r, \dim(V) = m + s$.

Lemma 2.15. $w_1, ..., w_m, u_1, ..., u_r, v_1, ..., v_s$ is a basis for U + V.

Assuming the above Lemma, the theorem is proved, for then $\dim(U + V) = m + r + s$, so

$$\dim(\mathbf{U}) + \dim(\mathbf{V}) = (\mathbf{m} + \mathbf{r}) + (\mathbf{m} + \mathbf{s}) = \mathbf{m} + (\mathbf{m} + \mathbf{r} + \mathbf{s}) = \dim(\mathbf{U} \cap \mathbf{V}) + \dim(\mathbf{U} + \mathbf{V}).$$

Proof of Lemma 2.15. First, spanning. So let $z \in U + V$. Then z = u + v for some $u \in U$, $v \in V$. We can write

$$u = a_1w_1 + \ldots + a_mw_m + b_1u_1 + \ldots + b_ru_r, v = c_1w_1 + \ldots + c_mw_m + d_1v_1 + \ldots + d_sv_s,$$

so

$$z = u + v = (a_1 + c_1)w_1 + \ldots + (a_m + c_m)w_m + b_1u_1 + \ldots + b_ru_r + d_1v_1 + \ldots + d_sv_s,$$

so our set spans U + V.

Next, linear independence. So suppose

(2)
$$a_1w_1 + \ldots + a_mw_m + b_1u_1 + \ldots + b_ru_r + c_1v_1 + \ldots + c_cv_s = 0$$

and put $w := \sum a_i w_i$, $u = \sum b_i u_i$ and $v = \sum c_i v_i$. Now $v = -w - u \in U$, so $v \in U \cap V$, so $v = d_1 w_1 + \ldots + d_m w_m$ for some d_i . Thus:

$$\mathbf{v} = \mathbf{c}_1 \mathbf{v}_1 + \ldots + \mathbf{c}_s \mathbf{v}_s = \mathbf{d}_1 \mathbf{w}_1 + \ldots + \mathbf{d}_m \mathbf{w}_m.$$

Hence $c_1v_1 + \ldots + c_sv_s - d_1w_1 - \ldots - d_mw_m = 0$. But $w_1, \ldots, w_m, v_1, \ldots, v_s$ are linearly independent, so $c_1 = \ldots = c_s = d_1 = \ldots = d_m = 0$, so v = 0. Thus, (2) gives

 $a_1w_1+\ldots+a_mw_m+b_1u_1+\ldots+b_bu_r=0,$

so (as $w_1, \ldots, w_m, u_1, \ldots, u_r$ are linearly independent), $a_1 = \ldots = a_m = b_1 = \ldots = b_r = 0$ as required.

Corollary 2.16. *If* $V = V_1 \oplus V_2$ *, then* $dim(V) = dim(V_1) + dim(V_2)$ *.*

Proof. This follows from the theorem, as in this case $\dim(V_1 \cap V_2) = 0$.

Example 2.17. Inside $V = \mathbb{R}^4$, let

$$\begin{array}{rcl} V_1 &:= & \{(a,b,c,0):a,b,c\in \mathbb{R}\},\\ V_2 &:= & \{(a,a,b,b):a,b\in \mathbb{R}\}. \end{array}$$

Then $V_1 \cap V_2 := \{(a, a, 0, 0) : a \in \mathbb{R}\} := \{a(1, 1, 0, 0) : a \in \mathbb{R}\}$, so dim $(V_1 \cap V_2) = 1$. Clearly dim $(V_1) = 3$, and also dim $(V_2) = 2$ (V_2 has basis (1, 1, 0, 0), (0, 0, 1, 1)). Now

$$\dim(V_1 + V_2) = \dim(V_1) + \dim(V_2) - \dim(V_1 \cap V_2) = 3 + 2 - 1 = 4$$

Thus, $V_1 + V_2 = \mathbb{R}^4$ (but it is not a direct sum).

3. LINEAR TRANSFORMATIONS

Definition 3.1. Let U, V be vector spaces over K. A mapping $\alpha : U \rightarrow V$ is a *linear transformation* (or linear mapping) if

$$\alpha(\mathfrak{u}_1 + \mathfrak{u}_2) = \alpha(\mathfrak{u}_1) + \alpha(\mathfrak{u}_2)$$
 and
 $\alpha(\mathfrak{a}\mathfrak{u}) = \mathfrak{a}\alpha(\mathfrak{u})$

for all $u_1, u_2, u \in U, a \in K$.

Equivalently, α is a linear transformation if $\alpha(a_1u_1 + a_2u_2) = a_1\alpha(u_1) + a_2\alpha(u_2)$ for all $a_1, a_2 \in K$ and $u_1, u_2 \in U$.

Exercise 3.2. Show that under linear transformation:

(1) The null vector of U is mapped to the null vector of V.

 \square

(2) Any straight line in U is mapped to a straight line in V.

Theorem 3.3. Any linear transformation from \mathbb{R}^m to \mathbb{R}^n has the form

$$\alpha(\mathbf{x}_1,\ldots,\mathbf{x}_m)^{\mathsf{T}}=\mathsf{A}(\mathbf{x}_1,\ldots,\mathbf{x}_m)^{\mathsf{T}},$$

where $A = (a_{ij})$ is an $n \times m$ matrix. The basis vectors $(1, 0, ..., 0)^T, ..., (0, ..., 0, 1)^T$ map under α to the columns of A.

Remark 3.4. Here, to save on trees, we write $(x_1, ..., x_m)^T$ (the transpose) for the column vector $\begin{pmatrix} x_1 \\ \vdots \\ x_m \end{pmatrix}$.

Example 3.5. Consider linear mappings $\mathbb{R}^2 \to \mathbb{R}^2$. As examples, we have the identity transformation (matrix $\begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}$), expansions (with matrix $\begin{pmatrix} a & 0 \\ 0 & a \end{pmatrix}$), rotations (matrix $\begin{pmatrix} \cos \theta & -\sin \theta \\ \sin \theta & \cos \theta \end{pmatrix}$ for rotation anticlockwise by θ) and a reflection in the x-axis (matrix $\begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix}$).

Before proving the theorem, we fix the following notation, used throughout the course. We let $e_1 := (1, 0, ..., 0)^T$, $e_2 := (0, 1, 0, ..., 0)^T$, ..., $e_m := (0, ..., 0, 1)^T$ be the *standard basis* of \mathbb{R}^m .

Proof of Theorem **3.3***.* By linearity of α *,*

$$\alpha(x_1,\ldots,x_m)^{\mathsf{T}} = \alpha(\sum_{j=1}^m x_j e_j) = \sum_{j=1}^m x_j \alpha(e_j) = \sum_{j=1}^m x_j f_j,$$

where $f_1 := \alpha(e_1), \ldots, f_m := \alpha(e_m)$. Write $f_j = (a_{1j}, \ldots, a_{nj})^T$, the jth column of A. Then $\alpha(x_1, \ldots, x_m)^T = \sum_{j=1}^m x_j f_j =$

$$= x_1 \begin{pmatrix} a_{11} \\ \vdots \\ a_{n1} \end{pmatrix} + x_2 \begin{pmatrix} a_{12} \\ \vdots \\ a_{n2} \end{pmatrix} + \ldots + x_m \begin{pmatrix} a_{1m} \\ \vdots \\ a_{nm} \end{pmatrix} = \begin{pmatrix} a_{11} & \ldots & a_{1m} \\ \vdots & \vdots & \vdots \\ a_{n1} & \ldots & a_{nm} \end{pmatrix} \begin{pmatrix} x_1 \\ \vdots \\ x_m \end{pmatrix},$$

as claimed in the theorem.

Remark 3.6. If we choose any $f_1, \ldots f_m \in \mathbb{R}^n$, then there is a *unique* linear mapping β : $\mathbb{R}^m \to \mathbb{R}^n$ with $\beta(e_i) = f_i$ for each $j = 1, \ldots, m$. For we must have $\beta(\sum_{j=1}^m x_j e_j) = \sum_{j=1}^m x_j f_j$, by linearity.

Now we do something similar for linear transformations between *arbitrary* vector spaces.

Theorem 3.7. Let U, V be vector spaces over K of dimensions m, n respectively, and let u_1, \ldots, u_m be a basis of U, $v_1, \ldots v_n$ a basis of V. Let $\alpha : U \to V$ be a linear transformation. For each $j = 1, \ldots, m$, put $\alpha(u_j) = \sum_{i=1}^n a_{ij}v_i$, and let $A = (a_{ij})$, an $n \times m$ matrix. Then for any $c_1, \ldots, c_m \in K$, $\alpha(\sum_{j=1}^m c_j u_j) = \sum_{i=1}^n d_i v_i$, where $(d_1, \ldots, d_n)^T = A(c_1, \ldots, c_m)^T$.

Proof. We have

$$\alpha \left(\sum_{j=1}^{m} c_{j} u_{j} \right) = \sum_{j=1}^{m} c_{j} \alpha(u_{j}) \quad \text{(by linearity)}$$
$$= \sum_{j=1}^{m} c_{j} \sum_{i=1}^{n} a_{ij} v_{i}$$
$$= \sum_{j=1}^{m} \sum_{i=1}^{n} a_{ij} c_{j} v_{i}$$
$$= \sum_{i=1}^{n} \left(\sum_{j=1}^{m} a_{ij} c_{j} \right) v_{i}$$
$$= \sum_{i=1}^{n} d_{i} v_{i}.$$

Remark 3.8. The j-th column of A gives the coordinates of $\alpha(u_j)$ with respect to the basis v_1, \ldots, v_n of V.

Much of the aim of the rest of the module is to choose bases carefully so that the matrix for A (which depends on the choice of bases) has a nice shape.

Definition 3.9. Let $\alpha : U \to V$ be a linear transformation. Then the *image* of α is $\text{Im}(\alpha) := {\alpha(u) : u \in U}$, and the *kernel* of α is $\text{Ker}(\alpha) := {u \in U : \alpha(u) = 0}$. The kernel is also called the *null space* of α .

Theorem 3.10. Let $\alpha : U \to V$ be a linear transformation. Then $Im(\alpha)$ is a subspace of V and $Ker(\alpha)$ is a subspace of U.

Proof. First, observe that $\alpha(0) = \alpha(0 \cdot u) = 0 \cdot \alpha(u) = 0$, so $0_V \in \text{Im}(\alpha)$, and $0_U \in \text{Ker}(\alpha)$. To see $\text{Im}(\alpha)$ is a subspace: if $\nu_1, \nu_2 \in \text{Im}(\alpha)$ then there are $u_1, u_2 \in U$ with $\alpha(u_1) = \nu_1$,

 $\alpha(u_2)=\nu_2. \text{ Now if } a_1,a_2\in K \text{, then } a_1\nu_1+a_2\nu_2=\alpha(a_1u_1+a_2u_2)\in Im(\alpha).$

To see Ker(α) is a subspace, let $u_1, u_2 \in \text{Ker}(\alpha)$. Then $\alpha(a_1u_1 + a_2u_2) = a_1\alpha(u_1) + a_2\alpha(u_2) = a_1 \cdot 0 + a_2 \cdot 0 = 0$, so $a_1u_1 + a_2u_2 \in \text{Ker}(\alpha)$.

Example 3.11. Let $\alpha : \mathbb{R}^3 \to \mathbb{R}^3$ be the linear map

$$(\mathbf{x},\mathbf{y},\mathbf{z})^{\mathsf{T}}\mapsto (\mathbf{x}+\mathbf{y}-2\mathbf{z},\mathbf{x}-\mathbf{z},\mathbf{x}-\mathbf{y})^{\mathsf{T}}.$$

We shall find the matrix for α (with respect to the standard basis of \mathbb{R}^3) and then find bases for Im(α) and Ker(α). So we see how α acts on the standard basis. Now $\alpha(1,0,0)^T = (1,1,1)^T$, $\alpha(0,1,0)^T = (1,0,-1)^T$, and $\alpha(0,0,1)^T = (-2,-1,0)^T$. Hence the matrix for α is

$$\mathsf{A} = egin{pmatrix} 1 & 1 & -2 \ 1 & 0 & -1 \ 1 & -1 & 0 \end{pmatrix}.$$

Now Ker(α) consists of all solutions to the simultaneous equations x+y-2z = 0, x-z = 0 and x-y = 0. By row operations, this solution space is { $a(1,1,1) : a \in \mathbb{R}$ } so {(1,1,1)} is a

basis for Ker(α). Also, Im(α) is spanned by the columns of A. To get a basis, write these as rows (ie form A^T), and do row operations. We can reduce A^T to

$$\begin{pmatrix} 1 & 1 & 1 \\ 0 & -1 & -2 \\ 0 & 0 & 0 \end{pmatrix},$$

so the non-zero rows $\{(1, 1, 1), (0, -1, -2)\}$ form a basis for Im(α).

Definition 3.12. If $\alpha : U \to V$ is a linear transformation, then $r(\alpha) := \dim(\operatorname{Im}(\alpha))$ is the *rank* of α , and $n(\alpha) = \dim(\operatorname{Ker}(\alpha))$ is the *nullity* of α .

Theorem 3.13 (Rank and Nullity Theorem). Let $\alpha : U \to V$ be a linear transformation. Then

$$\mathbf{r}(\alpha) + \mathbf{n}(\alpha) = \dim(\mathbf{U}).$$

Proof. As usual, we shall assume dim(U) is finite.

Let v_1, \ldots, v_r be a basis for $\text{Im}(\alpha)$. Then there are $u_1, \ldots, u_r \in U$ with $\alpha(u_i) = v_i$ for each i. Also, choose a basis t_1, \ldots, t_k for $\text{Ker}(\alpha)$. The theorem follows from the following claim.

Lemma 3.14. $u_1, \ldots, u_r, t_1, \ldots, t_k$ is a basis for U.

Proof of Lemma. Spanning. If $u \in U$, then $\alpha(u) \in Im(\alpha)$, so we can write $\alpha(u) = a_1v_1 + ... + a_rv_r = a_1\alpha(u_1) + ... + a_r\alpha(u_r)$. Then

$$\alpha(\mathfrak{u}-(\mathfrak{a}_1\mathfrak{u}_1+\ldots+\mathfrak{a}_r\mathfrak{u}_r))=\alpha(\mathfrak{u})-(\mathfrak{a}_1\alpha(\mathfrak{u}_1)+\ldots+\mathfrak{a}_r\alpha(\mathfrak{u}_r))=0.$$

Hence $u - (a_1u_1 + \ldots + a_ru_r) \in \text{Ker}(\alpha)$, so has form $b_1t_1 + \ldots + b_kt_k$ for some scalars b_i . Then $u = a_1u_1 + \ldots + a_ru_r + b_1t_1 + \ldots + b_kt_k$.

Linear Independence. Suppose $a_1u_1 + \ldots + a_ru_r + b_1t_1 + \ldots + b_kt_k = 0$. Then

 $0 = \alpha(a_1u_1 + \ldots + a_ru_r + b_1t_1 + \ldots + b_kt_k) = a_1v_1 + \ldots + a_rv_r.$

Hence $a_1 = \ldots = a_r = 0$, as v_1, \ldots, v_r are linearly independent. Then $b_1t_1 + \ldots + b_kt_k = 0$. But t_1, \ldots, t_k are linearly independent, so also $b_1 = \ldots = b_k = 0$.

 \square

Remark 3.15. In Example 3.11, dim(U) = dim(\mathbb{R}^3) = 3, n(α) = 1, and r(α) = 2. The proof of Theorem 3.13 shows that $(1, 1, 1)^T$, $\alpha^{-1}(1, 1, 1)^T$, $\alpha^{-1}(0, -1, -2)^T$ is a basis for \mathbb{R}^3 .

3.1. **Discussion of Simultaneous Equations.** Suppose we have a system of n linear simultaneous equations

$$\begin{array}{rcl} a_{11}x_1 + \ldots + a_{1m}x_m & = & b_1, \\ a_{21}x_1 + \ldots + a_{2m}x_m & = & b_2 \\ & \ddots & \ddots & \ddots \\ a_{n1}x_1 + \ldots + a_{nm}x_m & = & b_n, \end{array}$$

so of the form Ax = b, where $A = (a_{ij})$ is an $n \times m$ matrix, $x = (x_1, \dots, x_m)^T$, and $b = (b_1 \dots, b_n)^T \in \mathbb{R}^n$.

Multiplication by A gives a linear map $\alpha : \mathbb{R}^m \to \mathbb{R}^n$. Now $r(\alpha) = n \Leftrightarrow Im(\alpha) = \mathbb{R}^n$, which holds if and only if we can solve the above system for *every* $b \in \mathbb{R}^n$.

Also, $n(\alpha) = 0$ if and only if the homogeneous system $Ax = (0, ..., 0)^T$ has the zero vector in \mathbb{R}^m as the only solution. This means also that if $\alpha(x) = \alpha(y) = b$, then $\alpha(x-y) = 0$, so x - y = 0; that is, $n(\alpha) = 0$ means that solutions to Ax = b, if they exist, are *unique*.

Now $r(\alpha) + n(\alpha) = m$, so if m < n, then $r(\alpha) < n$ and we can't always solve the system. If m > n, then $m > r(\alpha)$ so $n(\alpha) > 0$, wo we cannot solve *uniquely*. If m = n, then either $r(\alpha) = n$ and $n(\alpha) = 0$ (so for all $b \in \mathbb{R}^n$ there is a *unique* solution of Ax = b), or $r(\alpha) < n$ and $n(\alpha) > n$ (so for some b there is no solution, but if there is a solution there are infinitely many).

3.2. Composition of Mappings.

Lemma 3.16. Let U, V, W be vector spaces, and $\alpha : U \to V, \beta : V \to W$ be linear transformations. Then the function $\beta \alpha : U \to W$ defined by $\beta \alpha(u) = \beta(\alpha(u))$ is linear.

Proof. We have

$$\begin{split} (\beta \alpha)(a_1u_1 + a_2u_2) &= & \beta(\alpha(a_1u_1 + a_2u_2)) \\ &= & \beta(a_1\alpha(u_1) + a_2\alpha(u_2)) & \text{(as α is linear)} \\ &= & a_1\beta(\alpha(u_1)) + a_2\beta(\alpha(u_2)) & \text{(as β is linear)} \\ &= & a_1(\beta\alpha)(u_1) + a_2(\beta\alpha)(u_2). \end{split}$$

Theorem 3.17. Suppose that dim(U) = m, with basis u_1, \ldots, u_m , dim(V) = n, with basis v_1, \ldots, v_n , dim(W) = p, with basis w_1, \ldots, w_p . Suppose $\alpha : U \rightarrow V$ and $\beta : V \rightarrow W$ are linear transformations with matrices A, B with respect to these bases. Then $\beta \alpha$ has matrix BA with respect to the bases u_1, \ldots, u_m of U and w_1, \ldots, w_p of W.

Proof. By Theorem 3.3,

$$\begin{aligned} \alpha(u_j) &= \sum_{i=1}^n a_{ij} v_i \text{ for } j = 1, \dots, m \\ \beta(v_i) &= \sum_{k=1}^p b_{ki} w_k \text{ for } i = 1, \dots, n, \end{aligned}$$

where $A = (a_{ij})$ and $B = (b_{ij})$. Put $C = BA = (c_{ij})$. Now,

$$(\beta \alpha)(u_j) = \beta(\sum_{i=1}^n a_{ij}v_i)$$

=
$$\sum_{i=1}^n a_{ij}\beta(v_i)$$

=
$$\sum_{i=1}^n \sum_{k=1}^p b_{ki}a_{ij}w_k$$

=
$$\sum_{k=1}^p (\sum_{i=1}^n b_{ki}a_{ij})w_k$$

=
$$\sum_{k=1}^p c_{kj}w_k.$$

Thus, $\beta \alpha$ has matrix C.

3.3. **Inverses.** We adopt the following convention. If $\alpha : U \rightarrow U$, then we use the same basis for U on both sides.

Definition 3.18. If $\alpha : U \to U$ is linear, then α is *invertible* if α is 1-1 and onto. In this case, for each $b \in U$, the equation $\alpha(u) = b$ has a unique solution, which we write as $u = \alpha^{-1}(b)$. We call α^{-1} the *inverse* of α .

Proposition 3.19. *Suppose* α : $U \rightarrow U$ *is invertible. Then*

- (1) $r(\alpha) = \dim(U), n(\alpha) = 0;$
- (2) α^{-1} : U \rightarrow U is linear.
- (3) if α has matrix A with respect to basis u₁,..., u_n, the α⁻¹ has matrix A⁻¹ with respect to this basis.

Proof. First, since α is onto, Im(α) = U, so r(α) = dim(U). Since α (u) = b can be solved *uniquely*, ker(α) = {0}, so n(α) = 0.

To see that α^{-1} is linear, suppose $\alpha^{-1}(v_1) = u_1$ and $\alpha^{-1}(v_2) = u_2$, and let a_1, a_2 be scalars. Then $\alpha(a_1u_1 + a_2u_2) = a_1v_1 + a_2v_2$, so

$$\alpha^{-1}(a_1v_1 + a_2v_2) = a_1u_1 + a_2u_2 = a_1\alpha^{-1}(v_1) + a_2\alpha^{-1}(v_2).$$

Finally, if α^{-1} has matrix B, then since $\alpha \alpha^{-1} = \alpha^{-1} \alpha = id$ which has matrix I_n , by Theorem 3.17 AB = BA = I, so $B = A^{-1}$.

Theorem 3.20 (The AP = PB Theorem). Let α : U \rightarrow U be a linear transformation, and u_1, \ldots, u_n and v_1, \ldots, v_n be two bases of U. Suppose that α is respresented by the matrix A with respect to u_1, \ldots, u_n , and by B with respect to basis v_1, \ldots, v_n . Then there is a nonsingular matrix P such that AP = PB, and P = (p_{ij}) is given by $v_j = \sum_{i=1}^{n} p_{ij} u_i$ (for each j).

Remark 3.21. (1) The conclusion says $B = P^{-1}AP$.

(2) P is a 'change of basis' matrix. It's jth column is (p_{1j},..., p_{nj})^T – the coordinates of ν_j with respect to the basis u₁,..., u_n.

Proof of Theorem 3.20. As α is represented by B in the v-basis,

$$\begin{aligned} \alpha(\nu_j) &= \sum_{i=1}^n b_{ij}\nu_i \\ &= \sum_{i=1}^n b_{ij}\sum_{k=1}^n p_{ki}u_k \\ &= \sum_{k=1}^n \sum_{i=1}^n p_{ki}b_{ij}u_k \\ &= \sum_{k=1}^n (PB)_{kj}u_k. \end{aligned}$$

Also,

$$\begin{aligned} \alpha(v_j) &= \alpha(\sum_{i=1}^n p_{ij}u_i) \\ &= \sum_{i=1}^n p_{ij}\alpha(u_i) \\ &= \sum_{i=1}^n \sum_{k=1}^n p_{ij}a_{ki}u_k \\ &= \sum_{k=1}^n \sum_{i=1}^n a_{ki}p_{ij}u_k \\ &= \sum_{k=1}^n (AP)_{kj}u_k. \end{aligned}$$

as A represents α with respect to u-basis

Thus, $(PB)_{kj} = (AP)_{kj}$ for all k, j, so PB = AP.

Finally, we show P is invertible. Now the jth column of P is v_j written in the u-basis. Thus, working in the u-basis, multiplication by P takes each u_i to v_i . This is invertible, with the inverse transformation taking v_i to u_i .

Example 3.22. Consider $\alpha : \mathbb{R}^2 \to \mathbb{R}^2$ given by $\alpha \begin{pmatrix} x \\ y \end{pmatrix} = \begin{pmatrix} x \\ -y \end{pmatrix}$. The u_1, u_2 basis will be $u_1 = \begin{pmatrix} 1 \\ 0 \end{pmatrix}, u_2 = \begin{pmatrix} 0 \\ 2 \end{pmatrix}$, and the v_1, v_2 basis will be $v_1 = \begin{pmatrix} 1 \\ 1 \end{pmatrix}, v_2 = \begin{pmatrix} 1 \\ 2 \end{pmatrix}$. First, we find A, the matrix of α for the u-basis. Now

$$\begin{aligned} \alpha(\mathfrak{u}_1) &= \alpha \begin{pmatrix} 1 \\ 0 \end{pmatrix} = \begin{pmatrix} 1 \\ 0 \end{pmatrix} = 1 \cdot \mathfrak{u}_1 + 0 \cdot \mathfrak{u}_2, \quad \text{and} \\ \alpha(\mathfrak{u}_2) &= \alpha \begin{pmatrix} 0 \\ 2 \end{pmatrix} = \begin{pmatrix} 0 \\ -2 \end{pmatrix} = 0 \cdot \mathfrak{u}_1 - 1 \cdot \mathfrak{u}_2. \end{aligned}$$

Hence $A = \begin{pmatrix} 1 & 0 \\ 0 & -1 \end{pmatrix}$.

Next, we find B, the matrix of α for the v-basis. Now

$$\alpha(\nu_1) = \alpha \begin{pmatrix} 1 \\ 1 \end{pmatrix} = \begin{pmatrix} 1 \\ -1 \end{pmatrix} = 3\nu_1 - 2\nu_2 \quad \text{and}$$
$$\alpha(\nu_2) = \alpha \begin{pmatrix} 1 \\ 2 \end{pmatrix} = \begin{pmatrix} 1 \\ -2 \end{pmatrix} = 4\nu_1 - 3\nu_2.$$

Hence, $B = \begin{pmatrix} 3 & 4 \\ -2 & -3 \end{pmatrix}$.

To find P: its jth column is v_j written in the u-basis.

$$egin{array}{rcl}
u_1&=&\begin{pmatrix}1\\1\end{pmatrix}=1\cdot\begin{pmatrix}1\\0\end{pmatrix}+rac{1}{2}\cdot\begin{pmatrix}0\\2\end{pmatrix}, & ext{ and}\
u_2&=&\begin{pmatrix}1\\2\end{pmatrix}=1\cdot\begin{pmatrix}1\\0\end{pmatrix}+1\cdot\begin{pmatrix}0\\2\end{pmatrix}, & ext{ } \end{array}$$

so $P = \begin{pmatrix} 1 & 1 \\ \frac{1}{2} & 1 \end{pmatrix}$. Finally, a calculation shows AP = PB.

4. DIAGONALISATION OF MATRICES

From now on, $\alpha : V \to V$ is a linear transformation, where V is an n-dimensional real or complex vector space. If α is represented by the $n \times n$ matrix A with respect to one fixed basis, then with respect to any other basis it has matrix P⁻¹AP, for some non-singular P (by Theorem 3.20).

Definition 4.1. Two matrices A, B are said to be *similar* if there is an invertible matrix P such that $P^{-1}AP = B$; that is, if A, B represent the same linear transformation with respect to (different) bases.

Exercise 4.2. Show that the relation "matrices A and B are similar" is an equivalence.

To investigate similarity of matrices we need the following notions.

Definition 4.3. Let $\alpha : V \to V$ be linear. A scalar λ is said to be an *eigenvalue* of α if there is a **non-zero** vector $v \in V$ such that $\alpha(v) = \lambda v$. The vector v is called an *eigenvector* corresponding to the eigenvalue λ .

Remark 4.4. (1) λ is an eigenvalue of $\alpha \Leftrightarrow$ there is $\nu \neq 0$ such that $(\alpha - \lambda I)\nu = 0$ $\Leftrightarrow \quad \text{Ker}(\alpha - \lambda I) \neq \{0\}$ $\Leftrightarrow \quad n(\alpha - \lambda I) \neq 0$ $\Leftrightarrow \quad r(\alpha - \lambda I) \neq n$ by Rank and Nullity Theorem $\Leftrightarrow \quad \alpha - \lambda I$ is not invertible.

(2) Some α have no eigenvalues. For example, suppose $K = \mathbb{R}$, and $\begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}$ defines a linear transformation $\alpha : \mathbb{R}^2 \to \mathbb{R}^2$ (with respect to the standard basis). Now $\begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix} \begin{pmatrix} x \\ y \end{pmatrix} = \lambda \begin{pmatrix} x \\ y \end{pmatrix}$ means that $y = \lambda x$ and $-x = \lambda y$, so either x = y = 0, or $-yx = \lambda^2 yx$, so $\lambda^2 = -1$, which has no solutions in \mathbb{R} .



FIGURE 1. Eigenvalues and Eigenvectors

Theorem 4.5. Let $\alpha : V \to V$ be a linear transformation, where dim(V) = n. Then the eigenvalues of α are the roots of the characteristic equation det(tI - A) = 0, where A is any matrix representing α with respect to some basis. The equation is independent of the choice of basis, so can be written det $(tI - \alpha) = 0$. It is a polynomial in t of degree n.

Proof. First, we verify that the characteristic equation is independent of the choice of basis. So suppose that A represents α with respect to one basis, and B represents α with respect to another. Then there is a non-singular matrix P such that $B = P^{-1}AP$ (Theorem 3.20). Now

$$det(tI - B) = det(tI - P^{-1}AP)$$

=
$$det(P^{-1}(tI)P - P^{-1}AP)$$

=
$$det(P^{-1}(tI - A)P)$$

=
$$det(P^{-1}) det(tI - A) det(P)$$

=
$$det(tI - A),$$

the last step using that $det(P^{-1}) det(P) = det(P^{-1}P) = det(I) = 1$.

Now, λ is an eigenvalue of $\alpha \Leftrightarrow \alpha - \lambda I$ is not invertible

$$\Leftrightarrow \det(\mathsf{A} - \lambda \mathsf{I}) = 0 \Leftrightarrow \det(\lambda \mathsf{I} - \mathsf{A}) = 0.$$

This clearly is a polynomial of degree n in t, with leading term t^n .

Remark 4.6. There is some confusion as to whether det(tI - A) or det(A - tI) is the characteristic equation. I go for det(tI - A), so the coefficient of t^n is 1. But the two have the same roots. Indeed, $det(tI - A) = (-1)^n det(A - tI)$.

We often write $\chi(t)$ for the characteristic equation (also called the *characteristic polynomial*). By Theorem 4.5, this can be regarded equally as the characteristic equation of A, or of α (where A represents α with respect to some basis).

- (1) The linear transformation $\alpha : V \to V$ is *diagonalisable* if there is Definition 4.7. a basis of V such that α is represented by a diagonal matrix with respect to this basis.
 - (2) We denote by $Diag(\lambda_1, \dots, \lambda_n)$ the diagonal matrix whose diagonal entry (in the (i, i) position, for each i) is λ_i .

Remark 4.8. If α is represented by A with respect to some basis, then α is diagonalisable if and only if there is a non-singular matrix P such that $P^{-1}AP$ is diagonal, that is, A is similar to a diagonal matrix.

Theorem 4.9. Let $\alpha : V \to V$ be linear. Then α is diagonalisable if and only if there is a basis of V consisting of eigenvectors of α .

Proof. \leftarrow If v_1, \ldots, v_n is a basis of V consisting of eigenvectors of α , with $\alpha(v_i) = \lambda_i v_i$ for each j, then the matrix for α with respect to the basis v_1, \ldots, v_n is $\text{Diag}(\lambda_1, \ldots, \lambda_n)$.

 \Rightarrow Conversely, suppose α has matrix $\text{Diag}(\lambda_1, \ldots, \lambda_n)$ with respect to basis ν_1, \ldots, ν_n . Then $\alpha v_i = \lambda_i v_i$ for each i, so the v_i are eigenvectors.

Example 4.10. Let $A = \begin{pmatrix} 0 & 1 \\ -1 & 0 \end{pmatrix}$. First, consider it as giving a linear transformation α : $\mathbb{R}^2 \to \mathbb{R}^2$ (with respect to the standard basis). The characteristic equation is det $\begin{pmatrix} \lambda & -1 \\ 1 & \lambda \end{pmatrix} =$

 $\lambda^2 + 1$, which has no real roots. So A has no eigenvalues in \mathbb{R} , so no eigenvectors, so is not diagonalisable. Next, work over *complex* scalars, and consider A as giving a linear transformation α : $\mathbb{C}^2 \to \mathbb{C}^2$, again with respect to the standard basis. Now $\lambda^2 + 1 = (\lambda + i)(\lambda - i)$, with roots $i_{,} -i$. We find the eigenvectors.

For
$$\lambda = i$$
, the equation $(iI - A) \begin{pmatrix} x \\ y \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \end{pmatrix}$ is
 $\begin{pmatrix} i & -1 \\ 1 & i \end{pmatrix} \begin{pmatrix} x \\ y \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \end{pmatrix}$,

which gives the single equation ix -y = 0, with solution set $\{a(1, i) : a \in \mathbb{C}\}$, so (1, i) is an eigenvector.

For
$$\lambda = -i$$
, we solve $(-iI - A) \begin{pmatrix} x \\ y \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \end{pmatrix}$, that is
 $\begin{pmatrix} -i & -1 \\ 1 & -i \end{pmatrix} \begin{pmatrix} x \\ y \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \end{pmatrix}$

This gives the equation -ix - y = 0, with solutions $\{a(i, 1) : a \in \mathbb{C}\}$, so (i, 1) is an eigenvector.

The vectors (1, i) and (i, 1) are linearly independent, so form a basis of eigenvectors of α , so by Theorem 4.9, α is dagonalisable. Pur P = $\begin{pmatrix} 1 & i \\ i & 1 \end{pmatrix}$ (so the *columns*) of P are the eigenvectors, treated as column vectors). Then $P^{-1} = \frac{1}{2} \begin{pmatrix} 1 & -i \\ -i & 1 \end{pmatrix}$, and a calculation

shows that $P^{-1}AP = \begin{pmatrix} i & 0 \\ 0 & -i \end{pmatrix}$.

4.1. More on eigenvalues, eigenvectors. In the discussion which follows, we shall work over \mathbb{C} . The key fact is that *every* polynomial $b_0 + b_1 t + \ldots + b_m t^m$ (with $b_i \in \mathbb{C}$) factorises completely into linear factors over \mathbb{C} . So we shall for a while assume $K = \mathbb{C}$. This avoids problems like that in the last example (working over \mathbb{R}).

Definition 4.11. Let $\alpha : V \to V$ be linear, and λ be an eigenvalue of α (so a root of the characteristic polynomial $\chi(t)$). Then the *algebraic multiplicity* $m_a(\lambda)$ is defined to be its multiplicity as a root of $\chi(t)$. Its *geometric multiplicity* $m_g(\lambda)$ is dim(ker(tI- α)) = n(tI- α).

Given an $n \times n$ matrix A, we can regard A as giving a linear transformation $\mathbb{C}^n \to \mathbb{C}^n$, so can talk of the algebraic and geometric multiplicities of eigenvalues of A (in fact, these do not depend on the choice of basis).

Proposition 4.12. Let $\alpha : V \to V$ be linear, where V is an n-dimensional vector space.

(1) For each eigenvalue λ of α , we have

$$1 \leq \mathfrak{m}_{\mathfrak{g}}(\lambda) \leq \mathfrak{m}_{\mathfrak{a}}(\lambda) \leq \mathfrak{n}$$

(that is, 'geometric multiplicity \leq algebraic multiplicity').

(2) If $\lambda_1, \ldots, \lambda_r$ are the distinct eigenvalues of α then

$$\mathfrak{n}_{\mathfrak{a}}(\lambda_1) + \mathfrak{m}_{\mathfrak{a}}(\lambda_2) + \ldots + \mathfrak{m}_{\mathfrak{a}}(\lambda_r) = \mathfrak{n}.$$

- (3) If $\lambda_1, ..., \lambda_r$ are the distinct eigenvalues, and ν_i is an eigenvector of λ_i for each *i*, then $\{\nu_1, ..., \nu_r\}$ is linearly independent.
- *Proof.* (1) By the definition of *eigenvalue* $m_g(\lambda) \ge 1$; and as $\chi(t)$ has degree n, $m_a(\lambda) \le n$. So we must show that $m_g(\lambda) \le m_a(\lambda)$.

Suppose that $m_g(\lambda) = p$, take a basis v_1, \ldots, v_p for $Ker(\lambda I - \alpha)$, and extend it to a basis v_1, \ldots, v_n for V. With respect to the basis v_1, \ldots, v_n , α has matrix

$$\begin{pmatrix} \text{Diag}(\lambda,\ldots,\lambda) & S\\ Z & T \end{pmatrix}$$
,

where $\text{Diag}(\lambda, ..., \lambda)$ is $p \times p$, S is some $p \times (n-p)$ matrix, Z is the all-zero $(n-p) \times p$ matrix, and T is some $(n-p) \times (n-p)$ matrix. Now $\det(tI-\alpha) = (t-\lambda)^p \det(tI-T)$, so $\mathfrak{m}_a(\lambda) \ge p$.

- (2) We have $\det(tI-\alpha) = (t-\lambda_1)^{\mathfrak{m}_{\alpha}(\lambda_1)}(t-\lambda_2)^{\mathfrak{m}_{\alpha}(\lambda_2)}\dots(t-\lambda_r)^{\mathfrak{m}_{\alpha}(\lambda_r)}$, and $\det(tI-\alpha)$ has degree n.
- (3) We show by induction on m that if λ₁,..., λ_m are distinct eigenvalues with eigenvectors v₁,..., v_m, then v₁,..., v_m are linearly independent.

If m = 1, this is trivial, as eigenvectors are non-zero. So suppose

$$c_1v_1+\ldots+c_mv_m=0.$$

Multiplying (3) by λ_m ,

$$c_1\lambda_m\nu_1+\ldots+c_m\lambda_m\nu_m=0$$

Also, by applying α to (3), we have

 $c_1\lambda_1\nu_1+\ldots+c_m\lambda_m\nu_m=0.$

Hence, subtracting, $c_1(\lambda_1 - \lambda_m)v_1 + \ldots + c_{m-1}(\lambda_1 - \lambda_m)v_m = 0$. As v_1, \ldots, v_{m-1} are linearly independent (by induction), this gives

 $c_1(\lambda_1 - \lambda_m) = \ldots = c_{m-1}(\lambda_{m-1} - \lambda_m) = 0.$

As the λ_j are distinct, this forces $c_1 = \ldots = c_{m-1} = 0$, and from (3) we then also get $c_m = 0$.

Definition 4.13. If V_1, \ldots, V_r are subspaces of the vector space V, then V is a *direct sum* $V = V_1 \oplus \ldots \oplus V_r$ if each $v \in V$ is *uniquely* expressible as $v = v_1 + \ldots + v_r$ (for $v_i \in V_i$).

Remark 4.14. The last definition extends Definition 1.16 to the situation where there are more than two subspaces.

Let V_1, \ldots, V_r be subspaces of V, and suppose that whenever $v_1 \in V_1, \ldots, v_r \in V_r$ and $v_1 + \ldots + v_r = 0$, we have $v_1 = \ldots = v_r = 0$. Then $V_1 \oplus \ldots \oplus V_r$ is a direct sum, and its dimension is dim $(V_1) + \ldots + \dim(V_r)$. We omit the proof, but it follows easily by induction from Corollary 2.16.

Proposition 4.15. *The following are equivalent.*

- (1) α is diagonalisable;
- (2) $m_q(\lambda) = m_a(\lambda)$ for each eigenvalue λ ;
- (3) $V = \operatorname{Ker}(\lambda_1 I \alpha) \oplus \ldots \oplus \operatorname{Ker}(\lambda_r I \alpha).$

Proof. For each i = 1, ..., r, put $V_i := \text{Ker}(\lambda_i I - \alpha)$.

(i) \Rightarrow (ii) We use Theorem 4.9 and its proof. So suppose α is diagonalisable. This means that with respect to a basis of eigenvectors, say u_1, \ldots, u_n , α has matrix

 $A = Diag(\lambda_1, \ldots, \lambda_1, \lambda_2, \ldots, \lambda_2, \ldots, \lambda_r, \ldots, \lambda_r).$

For each j there are $\mathfrak{m}_{\mathfrak{a}}(\lambda_j)$ occurrences of λ_j in this matrix (take its characteristic polynomial). Hence, there are $\mathfrak{m}_{\mathfrak{a}}(\lambda_j)$ basis vectors from $\mathfrak{u}_1, \ldots, \mathfrak{u}_n$ which are eigenvectors of λ_j (cf. Proof of Theorem 4.9). Hence, $\mathfrak{m}_{\mathfrak{g}}(\lambda_j) \ge \mathfrak{m}_{\mathfrak{a}}(\lambda_j)$. Thus, by (i), $\mathfrak{m}_{\mathfrak{g}}(\lambda_j) = \mathfrak{m}_{\mathfrak{a}}(\lambda_j)$ for each j.

(ii) ⇒ (iii) Suppose (ii), i.e. that m_g(λ) = m_a(λ) for each λ. First, using (iii) and Remark 4.14, observe that the sum V₁+...+V_r is direct. Thus, it suffices to show that it equals V, i.e. that the dimensions sum to n = dim(V). Since dim(V_i) = m_g(λ_i), assuming (b) we have

$$\sum_{i=1}^{r} dim(V_i) = \sum_{i=1}^{r} \mathfrak{m}_g(\lambda_i) = \sum_{i=1}^{r} \mathfrak{m}_a(\lambda_i) = \mathfrak{n},$$

as required.

(iii) \Rightarrow (i) Suppose V = V₁ \oplus ... \oplus V_r. Then, by Remark 4.14, dim(V) = dim(V₁)+...+dim(V_r). Let B_i be a basis for V_i for each i. Observe that B_i consists of eigenvectors with eigenvalue λ_i . Then, using the definition of direct sum, B := B₁ \cup ... \cup B_r is linearly independent, and as it has size dim(V), B is a basis of V. Thus, V has a basis of eigenvectors of α , so by Theorem 4.9, α is diagonalisable. Example 4.16. Let $A = \begin{pmatrix} 2 & 0 & 0 \\ 0 & 2 & 1 \\ 0 & 0 & 3 \end{pmatrix}$ and $B = \begin{pmatrix} 2 & 1 & 0 \\ 0 & 2 & 0 \\ 0 & 0 & 3 \end{pmatrix}$. In both cases, $\mathfrak{m}_{\mathfrak{a}}(3) = 1$ and

 $m_{\alpha}(2) = 2$, since each have characteristic equation $(\lambda - 2)^2(\lambda - 3)$. Since $1 \leq m_g(3) \leq m_{\alpha}(3)$, we must in each case have $m_g(3) = 1$. Thus, in each case, the matrix is diagonalisable if and only if $m_g(2) = 2$.

For A, we find $2I - A = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 0 & 1 \\ 0 & 0 & 1 \end{pmatrix}$, and the row reduced for of this has one non-zero

row. That is, r(2I - A) = 1 and n(2I - A) = 2 (regarding these as linear transformations with respect to standard basis), so for A we find $m_g(2) = 2$, and so A is diagonalisable.

For B, on the other hand, $2I - B = \begin{pmatrix} 0 & 1 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & 1 \end{pmatrix}$ which has rank 2 and nullity 1, so

 $m_g(2) = 1$ for B, so B is not diagonalisable. To diagonalise A find a basis of eigenvectors

To diagonalise A, find a basis of eigenvectors. For $\lambda = 2$, $v_1 = (1,0,0)^{\mathsf{T}}$ and $v_2 = (0,1,0)^{\mathsf{T}}$ are eigenvectors. For $\lambda = 3$, $v_3 = (0,1,1)^{\mathsf{T}}$ is an eigenvector. Let P have v_1, v_2, v_3 as columns, so $\mathsf{P} = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & 1 \\ 0 & 0 & 1 \end{pmatrix}$. Then $\mathsf{P}^{-1} = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 1 & -1 \\ 0 & 0 & 1 \end{pmatrix}$, and $\mathsf{P}^{-1}\mathsf{A}\mathsf{P} = \begin{pmatrix} 2 & 0 & 0 \\ 0 & 2 & 0 \\ 0 & 0 & 3 \end{pmatrix}$, which is discussed with the reduced probability of the transformation of the transfor

which is diagonal with the eigenvalues 2,2,3 listed in the order corresponding to eigenvectors v_1 , v_2 , v_3 .

4.2. **Polynomials.** We shall continue to work with vector spaces over \mathbb{C} , to ensure that any polynomial factorises into linear factors. So V is an n-dimensional vector space over \mathbb{C} .

We have $\alpha : V \to V$. If A represents α with respect to the basis v_1, \ldots, v_n , then by Theorem 3.17, A^2 represents the composition α^2 , A^3 represents α^3 , and so on. Also, if A represents α , and B represents β , then A+B represents $\alpha+\beta$, where $\alpha+\beta$ is defined so that $(\alpha+\beta)(\nu) = \alpha(\nu)+\beta(\nu)$. More generally, given the polynomial $p(t) = b_0+b_1t+\ldots+b_mt^m$, we can form the linear transformation $p(\alpha) = b_0 + b_1\alpha + \ldots + b_m\alpha^m$, which has matrix $p(A) := b_0I + b_1A + \ldots + b_mA^m$ with respect to the basis ν_1, \ldots, ν_n .

- **Lemma 4.17.** (1) Let $\alpha : V \to V$ be a non-zero linear transformation. Then there is a polynomial p(t) of degree at most n^2 with $p(\alpha) = 0$ (that is, $p(\alpha)$ is the zero linear transformation, so it maps each vector to the zero vector).
 - (2) All polynomials p such that $p(\alpha) = 0$ are multiples of a unique polynomial of minimal degree which is monic, *i.e.* has leading coefficient 1.

Definition 4.18. The monic polynomial of minimal degree described in (ii) is called the *minimal polynomial* of α , and is usually denoted $\mu(t)$.

Proof of Lemma 4.17. (1) Let Lin(V, V) be the set of all linear transformations $V \to V$. Then Lin(V, V) is a vector space over \mathbb{C} (compare Example 1.1(e).) If we work over a fixed basis of V, then Lin(V, V) is essentially the same as $M_{n,n}(\mathbb{C})$ (just replace each linear transformation by the matrix representing it in our given basis.) Now the dimension of $M_{n,n}(\mathbb{C})$ is n^2 , since we may form a basis { $E_{i,j} : 1 \leq i, j \leq n$ },

where $E_{i,j}$ has a 1 in the (i, j)-entry and zeros elsewhere. Hence dim $(Lin(V, V)) = n^2$. It follows that the set I, α , α^2 , ..., α^{n^2} is linearly *dependent* (as it has size $n^2 + 1$). Hence, there are constants $c_0, c_1, \ldots, c_{n^2} \in \mathbb{C}$ such that $c_0 + c_1 \alpha + \ldots c_{n^2} \alpha^{n^2} = 0$. Thus, if $p(t) = c_0 + c_1 t + \ldots + c_{n^2} t^{n^2}$, then $p(\alpha) = 0$.

(2) Let $\mu(t)$ be of minimal degree and of form $\mu(t) = t^k + d_{k-1}t^{k-1} + \ldots + d_1t + d_0 - we can put it in this (monic) form by multiplying by a non-zero constant.$

If p(t) is another polynomial with $p(\alpha) = 0$, then by division for polynomials, we have $p(t) = q(t)\mu(t) + r(t)$ where $deg(r(t)) < deg(\mu(t))$ or r(t) = 0. (If you don't know this fact about division of polynomials, don't worry about it! The idea is that q(t) is the quotient, r(t) is the remainder.) Now $r(\alpha) = p(\alpha)-q(\alpha)\mu(\alpha) = 0$, so by the minimality of the degree of $\mu(t)$, r(t) = 0, so $p(t) = q(t)\mu(t)$ is a multiple of $\mu(t)$.

We now state the main theorem about the minimal polynomial, the Cayley-Hamilton Theorem. The proof of the Cayley-Hamilton Theorem is given in these notes but not in the lectures, and for the exam, you will be expected to know the statement but not the proof of Cayley-Hamilton, and you will be expected to know the proof of Corollary 4.20.

Theorem 4.19 (Cayley-Hamilton Theorem). *The characteristic polynomial* $\chi(t)$ *of* α *satisfies* $\chi(\alpha) = 0$.

Corollary 4.20. If $\mu(t)$ and $\chi(t)$ are respectively the minimal and characteristic polynomials of t, then $\mu(t)$ divides $\chi(t)$.

Proof. By the Cayley-Hamilton Theorem, $\chi(\alpha) = 0$, so by Lemma 2, $\chi(t)$ is a multiple of $\mu(t)$.

Before proving the Cayley-Hamilton Theorem, recall that if A is an $n \times n$ matrix then adj(A) is the transposed 'matrix of cofactors'. That is, $adj(A) = B^T$, where $B = (b_{ij})$ and b_{ij} is + or - the determinant of the $(n-1) \times (n-1)$ -matrix obtained from A by removing

the ith row and jth column. For example, if $A = \begin{pmatrix} 1 & 2 & 1 \\ 0 & 1 & 1 \\ 1 & 0 & 1 \end{pmatrix}$ then $B = \begin{pmatrix} 1 & 1 & -1 \\ -2 & 0 & 2 \\ 1 & -1 & 1 \end{pmatrix}$,

so $adj(A) = B^{\mathsf{T}} = \begin{pmatrix} 1 & -2 & 1 \\ 1 & 0 & -1 \\ -1 & 2 & 1 \end{pmatrix}.$

Now $A \cdot adj(A) = det(A) \cdot I_n$ (a fact about matrices; you may recall that $A^{-1} = adj(A)/det(A)$ if A is invertible).

Proof of Theorem 4.19. Let A be the matrix of α with respect to some basis. By the last remark,

$$(tI - \alpha) adj(tI - \alpha) = \chi(t)I.$$

Put $B := adj(tI - \alpha)$. Then B is a matrix such that each entry is a polynomial in t. So $B = B_0 + tB_1 + t^2B_2 + \ldots$, where each B_j is a matrix over \mathbb{C} . (For example, we could write

$$\begin{pmatrix} t^2 + t + 1 & t - 1 \\ t^2 & t \end{pmatrix} = \begin{pmatrix} 1 & -1 \\ 0 & 0 \end{pmatrix} + \begin{pmatrix} 1 & 1 \\ 0 & 1 \end{pmatrix} t + \begin{pmatrix} 1 & 0 \\ 1 & 0 \end{pmatrix} t^2$$

Thus,

$$(tI - A)(B_0 + tB_1 + t^2B_2 + ...) = \chi(t)I,$$

for any t. Now put t = A. Then the left hand side is zero, so $\chi(A)I = 0$. This forces $\chi(A) = 0$, so $\chi(\alpha) = 0$.

Corollary 4.21. Suppose $\chi(t) = \prod_{i=1}^{m} (t - \lambda_i)^{n_i}$ is the characteristic polynomial of α . Then $\mu(t)$ has the form $\prod_{i=1}^{m} (t - \lambda_i)^{p_i}$ where $1 \leq p_i \leq n_i$ for each i.

Proof. Since $\mu(t)$ divides $\chi(t)$, we just need to check that each λ_i is a root of $\mu(t)$, i.e. that $1 \leq p_i$. So let λ be an eigenvalue of α , with eigenvector ν . Then $\alpha \nu = \lambda \nu$, so $\alpha^2 \nu = \lambda^2 \nu$, $\alpha^3(\nu) = \lambda^3(\nu)$, etc. Thus, $\mu(\alpha)\nu = \mu(\lambda)\nu$. Since $\mu(\alpha) = 0$, and $\nu \neq 0$, this forces $\mu(\lambda) = 0$, so λ is a root of $\mu(t)$.

Finally, we obtain our last criterion for diagonalisability of a linear transformation.

Proposition 4.22. Let $\alpha : V \to V$ be a linear transformation, where V is a vector space over \mathbb{C} . Then α is diagonalisable if and only if $\mu(t)$ is a product of distinct linear factors.

Proof. \Rightarrow Suppose that α is diagonalisable, represented by the matrix

$$A = \text{Diag}(\lambda_1, \ldots, \lambda_1, \lambda_2, \ldots, \lambda_2, \ldots, \lambda_r, \ldots, \lambda_r)$$

Then you can check that

$$(\lambda_1 I - A)(\lambda_2 I - A)\dots(\lambda_r I - A) = 0$$

so $\mu(t) = (t - \lambda_1)(t - \lambda_2) \dots (t - \lambda_r)$ (a product of *distinct* linear factors). \Leftarrow Suppose $\mu(t) = (t - \lambda_1)(t - \lambda_2) \dots (t - \lambda_r)$. Let $\mu_i(t) := \mu(t)/(t - \lambda_i)$. Now, expanding by partial fractions,

$$\frac{1}{\mu(t)} = \frac{1}{(t-\lambda_1)(t-\lambda_2)\dots(t-\lambda_r)} = \frac{c_1}{t-\lambda_1} + \frac{c_2}{t-\lambda_2} + \dots + \frac{c_r}{t-\lambda_r}.$$

Hence, multiplying by $\mu(t)$,

$$1 = \frac{\mu(t)}{\mu(t)} = c_1 \mu_1(t) + c_2 \mu_2(t) + \ldots + c_r \mu_r(t).$$

Let $v \in V$. Then, from the last equation,

(4)

$$\nu = c_1 \mu_1(\alpha) \nu + c_2 \mu_2(\alpha) \nu + \ldots + c_r \mu_r(\alpha) \nu.$$

It follows that for each i, $c_i \mu_i(\alpha) \nu \in \text{Ker}(\lambda_i I - \alpha)$. Indeed,

$$(\lambda_{i}I - \alpha)c_{i}\mu_{i}(\alpha)\nu = (\lambda_{i}I - \alpha)c_{i}\frac{\mu(\alpha)}{\lambda_{i}I - \alpha}\nu = c_{i}\mu(\alpha)\nu = 0.$$

Thus, each $c_i \mu_i(\alpha)$ is an eigenvalue of α , so by (4), the eigenvectors of α span V, and so α is diagonalisable by Theorem 4.9.

Example 4.23. Return to matrices A and B from the Example 4.16, which both have the characteristic polynomial $(\lambda - 2)^2(\lambda - 3)$. Since A is diagonalisable its minimal polynomial is $(\lambda - 2)(\lambda - 3)$, and since B is not diagonalisable its minimal polynomial coincides with the characteristic one. We also can conclude that A and B are not similar.

5. INNER PRODUCT SPACES

We now return to working with vector spaces over \mathbb{R} (so from now on, $K = \mathbb{R}$), but the theory below could be developed for vector spaces over \mathbb{C} .

Definition 5.1. Let V be a vector space over \mathbb{R} , and suppose that for each pair of vectors $u, v \in V$ there is defined a real number written $\langle u, v \rangle$ (sometimes denoted (u, v), or $u \cdot v$). This is called a (real) *inner product* on V if it satisfies the following, for $a, b \in \mathbb{R}$, and $u_1, u_2, u, v \in V$.

(1) (Linearity) $\langle au_1 + bu_2, v \rangle = a \langle u_1, v \rangle + b \langle u_2, v \rangle$.

(2) (Symmetry) $\langle u, v \rangle = \langle v, u \rangle$.

(3) (Positive Definiteness) $\langle u, u \rangle \ge 0$, and $\langle u, u \rangle = 0 \Leftrightarrow u = 0$.

The vector space V with $\langle -, - \rangle$ satisfying 1–3 is called an *inner product space*.

Remark 5.2. We can deduce from these axioms the following properties.

- (1) $\langle u, av_1 + bv_2 \rangle = \langle av_1 + bv_2, u \rangle = a \langle v_1, u \rangle + b \langle v_2, u \rangle = a \langle u, v_1 \rangle + b \langle u, v_2 \rangle$ that is, the inner product is linear also in the second variable.
- (2) For all $\nu \in V$, $\langle 0, \nu \rangle = \langle 0\nu, \nu \rangle = 0 \langle \nu, \nu \rangle = 0$, and also $\langle \nu, 0 \rangle = \langle 0, \nu \rangle = 0$.
- (3) We have the following generalisation of 1

$$\langle a_1u_1 + a_2u_2 + \ldots + a_nu_n, b_1v_1 + b_2v_2 + \ldots + b_mv_m \rangle = \sum_{i=1}^n \sum_{j=1}^m a_ib_j \langle u_i, v_j \rangle$$

This can be proved by induction. First use induction on n to get

$$\langle a_1u_1 + \ldots + a_nu_n, b_1v_1 \rangle = a_1b_1\langle u_1, v_1 \rangle + \ldots + a_nb_1\langle u_n, v_1 \rangle,$$

and then use induction on m.

- **Example 5.3.** (1) Take the usual scalar (or 'dot') product on $V = \mathbb{R}^n$. Here, if $x = (x_1, \ldots, x_n)$ and $y = (y_1, \ldots, y_n)$, then $\langle x, y \rangle = \sum_{i=1}^n x_i y_i$. This is 'Euclidean n-space'.
 - (2) Let V be the vector space of all continuous functions on the closed interval [a, b] (or restrict to the subspace of polynomial functions on this interval). Define

$$\langle f,g \rangle = \int_a^b f(t)g(t)dt.$$

This is an inner product, by the basic rules for integration.

(3) As a slight adjustment of (2), let *w*(t) be a continuous and strictly positive function on [a, b] (a 'weight function'), and let

$$\langle \mathbf{f}, \mathbf{g} \rangle = \int_{a}^{b} \mathbf{f}(\mathbf{t}) \mathbf{g}(\mathbf{t}) w(\mathbf{t}) d\mathbf{t}.$$

Theorem 5.4 (Cauchy-Schwarz Inequality). *If* V *is an inner product space, and* $u, v \in V$ *, then*

$$\langle \mathfrak{u}, \mathfrak{v} \rangle^2 \leqslant \langle \mathfrak{u}, \mathfrak{u} \rangle \langle \mathfrak{v}, \mathfrak{v} \rangle.$$

Proof. We can assume $\langle u, u \rangle \neq 0$, as otherwise u = 0 and both sides vanish. Let $t \in \mathbb{R}$. Then $\langle tu - v, tu - v \rangle \ge 0$, and (by linearity) equals

$$t^2 \langle u, u \rangle - 2t \langle u, v \rangle + \langle v, v \rangle.$$

Put $t = \frac{\langle u, v \rangle}{\langle u, u \rangle}$. Then we obtain

$$\frac{\langle \mathbf{u}, \mathbf{v} \rangle^2}{\langle \mathbf{u}, \mathbf{u} \rangle} - \frac{2 \langle \mathbf{u}, \mathbf{v} \rangle^2}{\langle \mathbf{u}, \mathbf{u} \rangle} + \langle \mathbf{v}, \mathbf{v} \rangle \geqslant 0.$$

Multiplying out, we get $\langle v, v \rangle \langle u, u \rangle \ge \langle u, v \rangle^2$.

Example 5.5. In \mathbb{R}^n , with the usual scalar product, the Cauchy-Schwarz Inequality says

$$\sum_{i=1}^{n} x_{i} y_{i} \leqslant \left(\sum x_{i}^{2}\right)^{\frac{1}{2}} \left(\sum y_{i}^{2}\right)^{\frac{1}{2}},$$

or $x \cdot y \leq |x| |y|$ for vectors in \mathbb{R}^n . This is true in \mathbb{R}^3 , as $x \cdot y = |x| |y| \cos \theta$, where θ is the angle between the vectors x and y.

Definition 5.6. If $v \in V$ then the *norm* of v, also called the *length* of v, is defined to be $\sqrt{\langle v, v \rangle}$, and denoted ||v||.

Theorem 5.7. Let V be an inner product space. Then the norm on V satisfies

- (1) $\|v\| \ge 0$, and $\|v\| = 0$ if and only if v = 0,
- (2) $||kv|| = |k| \cdot ||v||$ for any $k \in \mathbb{R}$,
- (3) $\|\mathbf{u} + \mathbf{v}\| \leq \|\mathbf{u}\| + \|\mathbf{v}\|$ (the triangle inequality).

Proof. (1) This is just an axiom for inner products.

(2) $\langle k\nu, k\nu \rangle = k^2 \langle \nu, \nu \rangle$. Now take square roots, noting both sides are positive.

(3) We have

$$\begin{aligned} \langle \mathbf{u} + \mathbf{v}, \mathbf{u} + \mathbf{v} \rangle &= \langle \mathbf{u}, \mathbf{u} \rangle + 2 \langle \mathbf{u}, \mathbf{v} \rangle + \langle \mathbf{v}, \mathbf{v} \rangle \\ &\leqslant \langle \mathbf{u}, \mathbf{u} \rangle + 2 \langle \mathbf{u}, \mathbf{u} \rangle^{\frac{1}{2}} \langle \mathbf{v}, \mathbf{v} \rangle^{\frac{1}{2}} + \langle \mathbf{v}, \mathbf{v} \rangle \quad \text{by Theorem 5.4} \\ &= \|\mathbf{u}\|^2 + 2 \|\mathbf{u}\| \cdot \|\mathbf{v}\| + \|\mathbf{v}\|^2 = (\|\mathbf{u}\| + \|\mathbf{v}\|)^2. \end{aligned}$$

Now take square roots.

Definition 5.8. We say that vectors $u, v \in V$ (a real inner product space) are *orthogonal*, written $u \perp v$, if $\langle u, v \rangle = 0$.

Theorem 5.9 (Pythagoras). *If* $u \perp v$ *then* $||u + v||^2 = ||u||^2 + ||v||^2$.

Proof. $\langle \mathbf{u} + \mathbf{v}, \mathbf{u} + \mathbf{v} \rangle = \langle \mathbf{u}, \mathbf{u} \rangle + 2 \langle \mathbf{u}, \mathbf{v} \rangle + \langle \mathbf{v}, \mathbf{v} \rangle = \|\mathbf{u}\|^2 + \|\mathbf{v}\|^2$.

Definition 5.10. A set $\{e_1, \ldots, e_m\}$ in an inner product space V is said to be *orthonormal* if $||e_i|| = 1$ for each i and $\langle e_i, e_j \rangle = 0$ for all distinct i, j. An *orthonormal basis* for a subspace W of V is a basis of W which is an orthonormal set.

Proposition 5.11. Suppose that e_1, \ldots, e_m is an orthonormal set.

- (1) If $v = \sum_{i=1}^{m} a_i e_i$ and $w = \sum_{i=1}^{m} b_i e_i$, then $\langle v, w \rangle = \sum_{i=1}^{m} a_i b_i$.
- (2) $\left\|\sum_{i=1}^{m} a_i e_i\right\|^2 = \left(\sum_{i=1}^{m} a_i^2\right), \text{ so } e_1, \dots, e_m \text{ is linearly independent;}$
- (3) If $v = a_1e_1 + \ldots + a_me_m$, then $a_i = \langle v, e_i \rangle$ for each i.

Proof. (1) Apply Remark 3 after Definition 5.1.

(2)

$$\langle \sum_{i=1}^m a_i e_i, \sum_{j=1}^m a_j e_j \rangle = \sum_{i=1}^m \sum_{j=1}^m a_i a_j \langle e_i, e_j \rangle = \sum_{i=1}^m a_i^2.$$

In particular, is $\sum_{i=1}^{m} a_i e_i = 0$, then $\sum_{i=1}^{m} a_i^2 = 0$, which implies that $a_i = 0$ for all i (as squares are non-negative).

(3) Observe that $\langle v, e_i \rangle = \langle a_1 e_1 + \ldots + a_m e_m, e_i \rangle = a_i$.

Example 5.12. (1) In \mathbb{R}^n , the standard basis is an orthonormal basis.

(2) There are many other orthonormal bases in \mathbb{R}^n , for example, in \mathbb{R}^3

$$u_1 = \left(\frac{1}{\sqrt{3}}, \frac{1}{\sqrt{3}}, \frac{1}{\sqrt{3}}\right),$$

$$u_2 = \left(\frac{1}{\sqrt{2}}, \frac{-1}{\sqrt{2}}, 0\right),$$

$$u_3 = \left(\frac{1}{\sqrt{6}}, \frac{1}{\sqrt{6}}, \frac{-2}{\sqrt{6}}\right)$$

is an orthonormal basis.

(3) Here is a set of orthonormal functions on the interval $[0, 2\pi]$, with respect to the inner product $\langle f, g \rangle = \int_0^{2\pi} f(t)g(t)dt$ (see Example 5.3):

$$\frac{1}{\sqrt{2\pi}}, \ \frac{\cos t}{\sqrt{\pi}}, \ \frac{\sin t}{\sqrt{\pi}}.$$

5.1. **Gram–Schmidt Orthogonalisation Process.** This is a procedure for finding an orthonormal basis. Let v_1, \ldots, v_m be a linearly independent set of vectors in the real inner product space V. We construct an orthonormal set e_1, \ldots, e_m such that for each $j = 1, \ldots, m$,

 $\operatorname{span}(e_1,\ldots,e_i) = \operatorname{span}(v_1,\ldots,v_i).$

The construction is as follows. First, put $e_1 = v_1 / ||v_1||$ (that is, we *normalise* v_1 , to ensure length 1).

To find each subsequent e_j , we first subtract off the components of v_j in the directions of the previous e_k , and then normalise.

More formally, put $w_2 = v_2 - \langle v_2, e_1 \rangle e_1$. Now

$$\langle w_2, e_1 \rangle = \langle v_2, e_1 \rangle - \langle v_2, e_1 \rangle \langle e_1, e_1 \rangle = 0,$$

and span $(w_2, e_1) = \text{span}(v_2, e_1) = \text{span}(v_2, v_1)$, so $w_2 \neq 0$. Now, put $e_2 = w_2 / ||w_2||$ (i.e., normalise w_2).

In general, let

$$w_j = v_j - \sum_{k=1}^{j-1} \langle v_j, e_k \rangle e_k.$$

Then, for l < j,

$$\langle w_{j}, e_{l} \rangle = \langle v_{j}, e_{l} \rangle - \sum_{k=1}^{j-1} \langle v_{j}, e_{k} \rangle \langle e_{k}, e_{l} \rangle = 0$$
, and

$$\operatorname{span}(w_j, e_1, \ldots, e_{j-1}) = \operatorname{span}(v_j, e_1, \ldots, e_{j-1}) = \operatorname{span}(v_1, \ldots, v_j)$$

so $w_j \neq 0$. Now put $e_j = w_j / ||w_j||$.

Example 5.13. The plane x + y + z = 0 in \mathbb{R}^3 has basis $v_1 = (1, 0, -1)$, $v_2 = (0, 1, -1)$. We find an orthonormal basis.

First, $e_1 = (1, 0, -1) / ||(1, 0, -1)|| = (\frac{1}{\sqrt{2}}, 0, \frac{-1}{\sqrt{2}}).$ Next,

$$w_2 = v_2 - \langle v_2, e_1 \rangle e_1 = (0, 1, -1) - \frac{1}{\sqrt{2}} \left(\frac{1}{\sqrt{2}}, 0, \frac{-1}{\sqrt{2}} \right) = \left(\frac{-1}{2}, 1, \frac{-1}{2} \right).$$

Hence,

$$e_2 = w_2 / ||w_2|| = \frac{\sqrt{2}}{\sqrt{3}} \left(\frac{-1}{2}, 1, \frac{-1}{2}\right) = \left(\frac{-1}{\sqrt{6}}, \frac{2}{\sqrt{6}}, \frac{-1}{\sqrt{6}}\right).$$

We'll now extend e_1 , e_2 to an orthonormal basis of \mathbb{R}^3 . Put $v_3 = (0, 0, 1)$, so v_1 , v_2 , v_3 is a basis of \mathbb{R}^3 . Then

$$w_{3} = v_{3} - \langle v_{3}, e_{1} \rangle e_{1} - \langle v_{3}, e_{2} \rangle e_{2}$$

$$= (0, 0, 1) - \frac{1}{\sqrt{2}} \left(\frac{1}{\sqrt{2}}, 0, \frac{-1}{\sqrt{2}} \right) + \frac{1}{\sqrt{6}} \left(\frac{-1}{\sqrt{6}}, \frac{2}{\sqrt{6}}, \frac{-1}{\sqrt{6}} \right)$$

$$= (0, 0, 1) + \left(\frac{1}{2}, 0, \frac{-1}{2} \right) + \left(\frac{-1}{6}, \frac{2}{6}, \frac{-1}{6} \right)$$

$$= \left(\frac{2}{6}, \frac{2}{6}, \frac{2}{6} \right) = \left(\frac{1}{3}, \frac{1}{3}, \frac{1}{3} \right).$$

Hence, $e_3 = \sqrt{3}w_3 = (\frac{1}{\sqrt{3}}, \frac{1}{\sqrt{3}}, \frac{1}{\sqrt{3}})$. Geometrically, e_3 is a unit vector perpendicular to the plane x + y + z = 0.

5.2. **Orthogonal Complements.** Recall from Definition 1.19: If $W_1 \subset V$ are vector spaces, then W_2 is a *complement* of W_1 if and only if $V = W_1 \oplus W_2$ (direct sum), or equivalently, $W_1 + W_2 = V$ and $W_1 \cap W_2 = \{0\}$, or equivalently again, each vector $v \in V$ is uniquely expressible as $v = w_1 + w_2$, for $w_i \in W_i$.

In general, complements to a subspace are not unique, see Remark 1.20. However, given an inner product we have a particularly nice complement.

Definition 5.14. Let V be an inner product space, and W be a subspace of V. The *orthogonal complement* W^{\perp} of W is defined by

$$W^{\perp} := \{ v \in V : \langle v, w \rangle = 0 \text{ for all } w \in W \}.$$

Theorem 5.15. Let W be a subspace of a finite-dimensional inner product space V. Then

- (1) W^{\perp} is a subspace of V;
- (2) $V = W \oplus W^{\perp}$, so dim $V = \dim W + \dim(W^{\perp})$;
- (3) $(W^{\perp})^{\perp} = W.$

Proof. (1) Clearly, $0 \in W^{\perp}$. If $v_1, v_2 \in W^{\perp}$ and $a_1, a_2 \in \mathbb{R}$, and $w \in W$, then

$$a_1v_1 + a_2v_2, w \rangle = a_1 \langle v_1, w \rangle + a_2 \langle v_2, w \rangle = 0 + 0 = 0.$$

Hence $a_1v_1 + a_2v_2 \in W^{\perp}$, so W^{\perp} is a subspace.

(2) Take a basis w₁,..., w_m for W, and extend it to a basis w₁,..., w_m, v_{m+1},..., v_n of V. Now apply Gram–Schmidt: We get an orthonormal basis e₁,..., e_n of V, such that e₁,..., e_m is an orthonormal basis of W. It suffices to prove the following claim, which yields that V = W + W[⊥] and W ∩ W[⊥] = {0}.

Lemma 5.16. $W^{\perp} = \text{span}(e_{m+1}, \dots, e_n).$

Proof of the Lemma 5.16. First, $W^{\perp} \subseteq \text{span}(e_{m+1}, \ldots, e_n)$. For let $\nu \in W^{\perp}$, and suppose $\nu = \sum_{i=1}^{n} a_i e_i$ (as the e_i are a basis of V). Then $a_j = \langle \nu, e_j \rangle = 0$ for $j \leq m$, since $e_j \in W$ and $\nu \in W^{\perp}$. Hence, $\nu = a_{m+1}e_{m+1} + \ldots + a_ne_n \in \text{span}(e_{m+1}, \ldots, e_n)$.

Also span(e_{m+1}, \ldots, e_n) $\subseteq W^{\perp}$. For suppose that $v \in \text{span}(e_{m+1}, \ldots, e_n)$ and $w \in W$, say $v = \sum_{i=m+1}^{n} a_i e_i$ and $w \in \sum_{i=1}^{m} b_i e_i$. Then

$$\langle v, w \rangle = \langle \mathfrak{a}_{\mathfrak{m}+1} \mathfrak{e}_{\mathfrak{m}+1} + \ldots + \mathfrak{a}_{\mathfrak{n}} \mathfrak{e}_{\mathfrak{n}}, \mathfrak{b}_1 \mathfrak{e}_1 + \ldots + \mathfrak{b}_{\mathfrak{m}} \mathfrak{e}_{\mathfrak{m}} \rangle = 0.$$

(3) First, we show $W \subseteq (W^{\perp})^{\perp}$. To see this, let $w \in W$ and $v \in W^{\perp}$. Then $\langle w, v \rangle = 0$, so $w \in (W^{\perp})^{\perp}$.

Thus, it remains to show that $\dim(W) = \dim(W^{\perp})$. But

$$\dim(W) + \dim(W^{\perp}) = n = \dim(W^{\perp})^{\perp} + \dim(W^{\perp})$$

(as $V = W^{\perp} + (W^{\perp})^{\perp}$ by 2), so dim $(W) = \dim(W^{\perp})^{\perp}$. Hence $W = (W^{\perp})^{\perp}$.

As a brief example, recall Example 5.13. There, $V = \mathbb{R}^3$, and W is the plane x+y+z = 0. Now

$$W^{\perp} = \operatorname{span}(e_3) = \operatorname{span}\left(\frac{1}{\sqrt{3}}, \frac{1}{\sqrt{3}}, \frac{1}{\sqrt{3}}\right).$$

6. REAL SYMMETRIC MATRICES

We show here that any real symmetric matrix can be diagonalised by an "orthogonal" matrix P. This gives applications to quadratic forms in the final section.

First, a quick observation about inner products. Let $\langle \cdot, \cdot \rangle$ be an inner product on the real vector space V, with an orthonormal basis e_1, \ldots, e_n for V. We shall write vectors $v \in V$ as column vectors with respect to the e_i basis (but for reasons of space, these will be denoted as transposes of row vectors). So if $v = \sum_{i=1}^{n} a_i e_i$, we write $v = (a_1, \ldots, a_n)^T$. Now observe that $\langle v, w \rangle = w^T v$ (a matrix product of a row vector by a column vector). For if $v = a_1e_1 + \ldots + a_ne_n$ and $w = b_1e_1 + \ldots + b_ne_n$, then

$$\langle \mathbf{v}, \mathbf{w} \rangle = \sum_{i=1}^{n} a_i b_i = (b_1, \dots, b_n) \begin{pmatrix} a_1 \\ \vdots \\ a_n \end{pmatrix} = \mathbf{w}^{\mathsf{T}} \mathbf{v}.$$

Recall that a matrix A is *symmetric* if $A = A^T$, that is, if $A = (a_{ij})$, then $a_{ij} = a_{ji}$ for all i, j.

Recall also that if matrix A represents the linear transformation $\alpha : V \to V$ with respect to the basis e_1, \ldots, e_n , then $\alpha(\sum_{j=1}^n x_j e_j) = \sum_{i=1}^n y_i e_i$, where $A(x_1, \ldots, x_n)^T = (y_1, \ldots, y_n)^T$, so $y_i = \sum_{j=1}^n a_{ij} x_j$ for each $i = 1, \ldots, n$.

 \square

 \square

Definition 6.1. The linear transformation $\alpha : V \to V$ is called *self-adjoint* if, for all $u, v \in V$, we have

$$\langle \alpha u, v \rangle = \langle u, \alpha v \rangle$$

Theorem 6.2. Let $\alpha : V \to V$ be linear, represented by A with respect to the orthonormal basis e_1, \ldots, e_n for V. Then A is symmetric if and only if α is self-adjoint.

Proof. \leftarrow Suppose $\langle \alpha u, v \rangle = \langle u, \alpha v \rangle$ for al $u, v \in V$. Put $u = e_i$ and $v = e_j$. Then, (with the subscripts i, j indicating which vector entry is 1),

$$\begin{split} \langle \alpha u, v \rangle &= v^{\mathsf{T}} \cdot \alpha u = (0, \dots, 0, 1_{j}, 0, \dots, 0) \mathsf{A}(0, \dots, 0, 1_{i}, 0, \dots 0)^{\mathsf{T}} \\ &= (0, \dots, 0, 1_{j}, 0, \dots, 0) \begin{pmatrix} a_{1i} \\ \vdots \\ a_{ni} \end{pmatrix} = a_{ji}. \end{split}$$

Also,

$$\langle \mathfrak{u}, \mathfrak{a} \mathfrak{v} \rangle = (\mathfrak{a}_{1j}, \ldots, \mathfrak{a}_{nj})(0, \ldots, 0, 1_i, 0, \ldots, 0)^{\mathsf{T}} = \mathfrak{a}_{ij}.$$

Hence, $a_{ij} = a_{ji}$.

⇒ Suppose A is symmetric. Let $v = (v_1, ..., v_n)^T$ and $u = (u_1, ..., u_n)^T$. Then $\langle \alpha u, v \rangle = v^T A u$, and $\langle u, \alpha v \rangle = (Av)^T u = v^T A^T u = v^T A u$ (as A is symmetric). Hence, $\langle \alpha u, v \rangle = \langle u, \alpha v \rangle$.

Theorem 6.3. Let $\alpha : V \to V$ be self-adjoint, represented by the real symmetric matrix A with respect to the orthonormal basis e_1, \ldots, e_n of V. Then

- (1) The eigenvalues of A (and hence α) are all real;
- (2) eigenvalues corresponding to distinct eigenvectors are orthogonal;
- (3) V has an orthonormal basis of eigenvectors of A.

In particular, α is diagonalisable.

Proof. (1) Suppose that λ is an eigenvalue of A, with eigenvector $x_1e_1 + \ldots + x_ne_n$. Put $x := (x_1, \ldots, x_n)^T$. Write $\overline{\lambda}$ for the complex conjugate of λ , and work over \mathbb{C} . (So if $\lambda = a + bi$, then $\overline{\lambda} = a - bi$.) Also, for the vector x, write $\overline{x} = (\overline{x}_1, \ldots, \overline{x}_n)^T$. Now $Ax = \lambda x$, so $A\overline{x} = \overline{\lambda x} = \overline{\lambda}\overline{x}$. Taking transposes, it follows that $\overline{x}^T A^T = \overline{\lambda}\overline{x}^T$.

Hence, multiplying on the right by the $n \times 1$ matrix x,

$$\bar{\mathbf{x}}^{\mathsf{T}} \mathbf{A}^{\mathsf{T}} \mathbf{x} = \bar{\lambda} \bar{\mathbf{x}}^{\mathsf{T}} \mathbf{x}.$$

Also,

$$\lambda \bar{\mathbf{x}}^{\mathsf{T}} \mathbf{x} = \bar{\mathbf{x}}^{\mathsf{T}} (\lambda \mathbf{x}) = \bar{\mathbf{x}}^{\mathsf{T}} (A \mathbf{x}) = \bar{\mathbf{x}}^{\mathsf{T}} A^{\mathsf{T}} \mathbf{x},$$

the last step as $A = A^T$. Thus, $\bar{\lambda}\bar{x}^T x = \lambda \bar{x}^T x$. As $x \neq 0$, it follows that $\lambda = \bar{\lambda}$, that is, λ is real.

(2) Suppose $\lambda \neq \mu$ are eigenvalues, with eigenvectors ν , w respectively. So $\alpha x = \lambda x$ and $\alpha y = \mu y$. Then

$$\lambda \langle x, y \rangle = \langle \lambda x, y \rangle = \langle \alpha x, y \rangle = \langle x, \alpha y \rangle = \langle x, \mu y \rangle = \mu \langle x, y \rangle.$$

As $\lambda \neq \mu$, it follows that $\langle x, y \rangle = 0$.

(3) We use induction on $n = \dim V$. If n = 1 the result is obvious.

By 1, there is a real eigenvalue λ_1 of α , with eigenvector x_1 . We may suppose that $||x_1|| = 1$ (normalise). Put $W := (\operatorname{span}(x_1))^{\perp} < V$. Then $\dim(W) = n - 1$, by Theorem 5.15(ii).

We claim that $\alpha(W) \subseteq W$. To see this, observe that if $w \in W$, then

$$\langle \alpha w, x_1 \rangle = \langle w, \alpha x_1 \rangle = \langle w, \lambda_1 x_1 \rangle = \lambda_1 \langle w, x_1 \rangle = 0,$$

so $\alpha w \in W$.

It follows that α induces a self-adjoint mapping $W \rightarrow W$. Since dim(W) = n - 1, there is an orthonormal basis w_2, \ldots, w_n of eigenvectors of α in W. Now x_1, w_2, \ldots, w_n is an orthonormal basis of eigenvectors of α in V.

Definition 6.4. An $n \times n$ matrix P is *orthogonal* if $P^T = P^{-1}$, that is $P^T P = I$.

The above definition esentially says that P is orthogonal if and only if its columns form an orthonormal basis of \mathbb{R}^n . In the 2 × 2 case, an example is $\begin{pmatrix} \cos \theta & -\sin \theta \\ \sin \theta & \cos \theta \end{pmatrix}$.

Lemma 6.5. Let $\alpha : V \to V$ be represented by a matrix A with respect to some orthonormal basis e_1, \ldots, e_n of V. Then if we change to another orthonormal basis f_1, \ldots, f_n of V, the matrix for α in this basis has form $P^{-1}AP$, where P is orthogonal.

Proof. We have $f_j = \sum_{i=1}^{n} p_{ij} e_i$, where $P = (p_{ij})$ (see Theorem 3.20). Now

$$\langle f_j, f_k \rangle = \langle \sum_{i=1}^n p_{ij} e_i, \sum_{i=1}^n p_{ik} e_i \rangle = \sum_{i=1}^n p_{ij} p_{ik}.$$

Hence, $\sum_{i=1}^{n} p_{ij} p_{ik}$ takes value 1 if j = k, and zero otherwise. This just says $P^T P = I$. \Box

The following corollary is really the point of this section. It says in particular than any real symmetric matrix is diagonalisable, but is much stronger than this (because P is orthogonal).

Corollary 6.6. If A is a real symmetric matrix, then there is an orthogonal matrix P with $P^{-1}AP = P^{T}AP$ equal to a diagonal matrix D.

Proof. Take e_1, \ldots, e_n as the standard basis of \mathbb{R}^n , and f_1, \ldots, f_n as an orthonormal basis of eigenvectors of A (given by Theorem 6.3). Then the matrix D of A with respect to f_1, \ldots, f_n is diagonal (as it is a basis of eigenvectors). Also, by Lemma 6.5, there is orthogonal P with $P^{-1}AP = D$.

Example 6.7. Let α : $\mathbb{R}^3 \to \mathbb{R}^3$ be given by the matrix $A = \begin{pmatrix} 0 & 2 & 0 \\ 2 & 2 & 2 \\ 0 & 2 & 0 \end{pmatrix}$ with respect

to the standard basis. Note that A is symmetric, so there is an orthonormal basis of \mathbb{R}^3 consisting of eigenvectors of A. First, calculating det(tI – A) along the first row, we find $\chi(t) = t(t(t-2)-4) + 2(-2t) = t(t^2 - 2t - 8) = t(t-4)(t+2)$, so the eigenvalues are 0,4,-2.

For $\lambda = 0$, solving

$$\begin{pmatrix} 0 & -2 & 0 \\ -2 & -2 & -2 \\ 0 & -2 & 0 \end{pmatrix} \begin{pmatrix} x \\ y \\ z \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \\ 0 \end{pmatrix}$$

gives eigenvector $(1,0,-1)^T$. For $\lambda = 4$, we find eigenvector $(1,2,1)^T$. For $\lambda = -2$, an eigenvector is $(1,-1,1)^T$. Normalising these, we find

$$\mathsf{P} = \begin{pmatrix} \frac{1}{\sqrt{2}} & \frac{1}{\sqrt{6}} & \frac{1}{\sqrt{3}} \\ 0 & \frac{2}{\sqrt{6}} & \frac{-1}{\sqrt{3}} \\ \frac{-1}{\sqrt{2}} & \frac{1}{\sqrt{6}} & \frac{1}{\sqrt{3}} \end{pmatrix}.$$

Then

$$\mathbf{P}^{-1}\mathbf{A}\mathbf{P} = \mathbf{P}^{\mathsf{T}}\mathbf{A}\mathbf{P} = \begin{pmatrix} 0 & 0 & 0 \\ 0 & 4 & 0 \\ 0 & 0 & -2 \end{pmatrix}.$$

7. QUADRATIC FORMS

Let V be an n-dimensional real vector space, with basis v_1, \ldots, v_n . So any vector $x \in V$ can be written in the form $x = x_1v_1 + \ldots + x_nv_n$.

Definition 7.1. In this setting, a *quadratic form* is a function $q : V \to \mathbb{R}$ of the form

$$q(x) = q(x_1v_1 + \ldots + x_nv_n) = \sum_{i=1}^n \sum_{j=1}^n c_{ij}x_ix_j,$$

where each $c_{ij} \in \mathbb{R}$.

- **Example 7.2.** (1) The standard norm of vector produced by the inner product $\langle x, x \rangle = x_1^2 + x_2^2 + \ldots + x_n^2$ in \mathbb{R}^n is a quadratic form. Moreover a norm from any inner product is a quadratic form as well.
 - (2) Working in the standard basis of \mathbb{R}^2 , with $x = (x_1, x_2)$, the function $q(x) = 3x_1^2 + 4x_1x_2 x_2^2$ is a quadratic form. We can write q(x) as a matrix product

$$(\mathbf{x}_1\mathbf{x}_2) \begin{pmatrix} 3 & 4 \\ 0 & -1 \end{pmatrix} \begin{pmatrix} \mathbf{x}_1 \\ \mathbf{x}_2 \end{pmatrix} = (\mathbf{x}_1\mathbf{x}_2) \begin{pmatrix} 3\mathbf{x}_1 + 4\mathbf{x}_2 \\ -\mathbf{x}_2 \end{pmatrix}.$$

It can also be written as

$$(\mathbf{x}_1\mathbf{x}_2) \begin{pmatrix} 3 & 2\\ 2 & -1 \end{pmatrix} \begin{pmatrix} \mathbf{x}_1\\ \mathbf{x}_2 \end{pmatrix}.$$

The second matrix expression above suggests there is a connection between quadratic forms and real symmetric matrices. We now formalise this.

Proposition 7.3. Working with respect to a fixed basis of V, any quadratic form q(x) can be written in the form

$$q(\mathbf{x}) = (\mathbf{x}_1 \dots \mathbf{x}_n) A \begin{pmatrix} \mathbf{x}_1 \\ \vdots \\ \mathbf{x}_n \end{pmatrix},$$

where A is a real symmetric matrix. Furthermore, such an expression is unique, and every $n \times n$ real symmetric matrix gives rise to a quadratic form on V.

Proof. We just prove the first assertion. Suppose our form is

$$q(x) = q(x_1\nu_1 + \ldots + x_n\nu_n) = \sum_{i=1}^n \sum_{j=1}^n c_{ij}x_ix_j.$$

Put $a_{ij} := \frac{c_{ij} + c_{ji}}{2}$. Then $A = (a_{ij})$ is symmetric, and

$$q(x) = \sum_{i=1}^{b} \sum_{j=1}^{n} a_{ij} x_i x_j = (x_1 \dots x_n) A (x_1 \dots x_n)^{\mathsf{T}}.$$

The last proposition says that, if we write $x = x_1v_1 + ... + x_nv_n$ as the column vector $(x_1, ..., x_n)^T$, then $q(x) = x^T A x$, where A is real symmetric. We would like to change the basis of V, to ensure that the matrix A is diagonal.

Proposition 7.4. Suppose that w_1, \ldots, w_n is another basis of V, with $w_j = \sum_{i=1}^n m_{ij} v_i$, for each j. Put $M = (m_{ij})$. Then

- (1) $\sum_{i=1}^{n} x_i v_i = \sum_{j=1}^{n} y_j w_j$, where $(x_1 \dots x_n)^T = M(y_1 \dots y_n)^T$.
- (2) In the w_i basis, the quadratic form $q(x) = x^T A x$ becomes

$$q(\mathbf{x}) = (\mathbf{y}_1 \dots \mathbf{y}_n)(\mathbf{M}^\mathsf{T} \mathbf{A} \mathbf{M}) \begin{pmatrix} \mathbf{y}_1 \\ \vdots \\ \mathbf{y}_n \end{pmatrix},$$

and the matrix M^TAM is symmetric. (Here, $(y_1, ..., y_n)$ is the sequence of coordinates of the vector x with respect to the w basis, that is, $x = y_1w_1 + ... + y_nw_n$.)

 \square

Proof. (1) We indeed have:

(2) Now
$$q(x) = x^{T}Ax = (My)^{T}A(My) = y^{T}(M^{T}AM)y.$$

Definition 7.5. We say that square matrices A, B are *congruent* if there is non-singular M with $B = M^{T}AM$.

- *Remark* 7.6. (1) In Proposition 7.4, the jth column of M is the tuple of coordinates of w_i written in the *v*-basis.
 - (2) By Corollary 6.6, any real symmetric matrix is congruent to a diagonal matrix.
 - (3) In the first year Geometry module you used orthogonal transformations of \mathbb{R}^2 (which are rotations) to diagonalise quadratic forms in order to obtain canonical equations of ellipses, parabolas and hyperbolas, cf. Example 7.13.

Theorem 7.7. For any quadratic form q on V, there is a choice of basis with redspect to which q is diagonal, that is, q is given as

$$q(y) = y' Dy = y' Diag(d_1, \dots, d_n)y$$

where $D = Diag(d_1, \dots, d_n)$ is a diagonal matrix. Equivalently, $q(y) = d_1y_1^2 + d_2y_2^2 + \dots + d_ny_n^2$.

Proof. We must show that if A is a real symmetric matrix, then there is a non-singular matrix M such that M^TAM is diagonal. By Corollary 6.6, there is an orthogonal matrix P such that $P^{-1}AP$ is diagonal, and as P is orthogonal, we have $P^{-1} = P^{T}$. \square

The proof of Theorem 7.7 gives a method of diagonalising a quadratic form q:

- (1) Find the real symmetric matrix A and its characteristic polynomial $\chi(t)$;
- (2) Factorise $\chi(t)$ to find the eigenvalues
- (3) Find a basis of eigenvectors, let M have these as columns, and calculate $M^{T}AM$.

This process is slow – for example, it may be difficult to factorise $\chi(t)$. Since we do not require that $M^{T} = M^{-1}$, there are much easier methods. The main point is that premultiplying by a non-singular matrix M^T has the effect of doing a sequence of row operations to A, and postmultiplying by M has the effect of doing the corresponding sequence of column operations. So we aim to diagonalise A by doing a sequence of row operations and the corresponding sequence of column operations. Here's an easy example.

Example 7.8. We diagonalise the quadratic form $q = x^2 + 4xy + 2yz$. (I have switched from having variables x_1, \ldots, x_n to variables x, y, z – you will see both conventions in past

papers.) The matrix here is $\begin{pmatrix} 1 & 2 & 0 \\ 2 & 0 & 1 \\ 0 & 1 & 0 \end{pmatrix}$. To diagonalise it, we first aim to make the 2's in the (1,2) and (2,1) entries into 0. So do $r'_2 = r_2 - 2r_1$ and then $c'_2 = c_2 - 2c_1$, to get

$$\begin{pmatrix} 1 & 2 & 0 \\ 2 & 0 & 1 \\ 0 & 1 & 0 \end{pmatrix} \sim \begin{pmatrix} 1 & 2 & 0 \\ 0 & -4 & 1 \\ 0 & 1 & 0 \end{pmatrix} \sim \begin{pmatrix} 1 & 0 & 0 \\ 0 & -4 & 1 \\ 0 & 1 & 0 \end{pmatrix}.$$

This is still symmetric. To get rid of the 1's in the (3, 2) and (2, 3) entries, do $r'_3 = r_3 + r_2/4$ and then $c'_3 = c_3 + c_2/4$. We get $\begin{pmatrix} 1 & 0 & 0 \\ 0 & -4 & 1 \\ 0 & 0 & \frac{1}{4} \end{pmatrix}$ and then $\begin{pmatrix} 1 & 0 & 0 \\ 0 & -4 & 0 \\ 0 & 0 & \frac{1}{4} \end{pmatrix}$. Thus, with respect to a different basis, the quadratic form can be written in the diagonal form $q = u^2 - 4v^2 + 1$ $\frac{1}{4}w^{2}$.

In general, the algorithm for diagonalising a quadratic form is slightly more complicated. For example, in the last example, how would we have got rid of the 2's in the (1,2) and (2,1) entries if the (1,1) entry had been 0? Here is a sketch of the general procedure. It ensures that there is no non-zero entry in the first row or column except possibly in the (1,1) entry. After we have achieved that, then we apply the same process to the $(n-1) \times (n-1)$ matrix obtained by deleting the first row and column.

- (1) If $a_{11} \neq 0$, then for each i > 1 do operations $r'_i = r_i + ar_1$ and then $c'_i = c_i + ac_1$ for appropriate a.
- (2) If $a_{11} = 0$ but $a_{ii} \neq 0$ for some i > 1, interchange rows r_1 and r_i , and then interchange columns c_1 and c_i , to get a matrix with a_{ii} in the (1,1) entry. Then apply (i) above.
- (3) Suppose all the diagonal entries a_{ii} are zero. Find some j with $a_{ij} \neq 0$, and then do $r'_i = r_i + r_j$ and then $c'_i = c_i + c_j$. The (i, i) entry is now non-zero, so proceed as in (ii).

Example 7.9. Consider the quadratic form q = 4yz (a function of x, y, z). This has matrix $\begin{pmatrix} 0 & 0 & 0 \end{pmatrix}$

 $\begin{pmatrix} 0 & 0 & 0 \\ 0 & 0 & 2 \\ 0 & 2 & 0 \end{pmatrix}$. There is no non-zero entry in first row or column, and no non-zero diagonal

entry, so use method (iii). Do $r'_2 = r_2 + r_3$ then $c'_2 = c_2 + c_3$ to get $\begin{pmatrix} 0 & 0 & 0 \\ 0 & 4 & 2 \\ 0 & 2 & 0 \end{pmatrix}$. Now do

$$r'_{3} = r_{3} - \frac{1}{2}r_{2}$$
 and then $c'_{3} = c_{3} - \frac{1}{2}c_{2}$, to get $\begin{pmatrix} 0 & 0 & 0 \\ 0 & 4 & 0 \\ 0 & 0 & -1 \end{pmatrix}$. Thus, a diagonalised form is $q = 4v^{2} - w^{2}$.

A given quadratic form can usually be diagonalised in many different ways, i.e., there is no *unique* answer. However, the following definition and theorem give a partial uniqueness.

Definition 7.10. Let $q(y) = d_1y_1^2 + ... + d_ny_n^2$ be a quadratic form in diagonal form. Let P be the number of strictly positive d_i , and N be the number of strictly negative d_i . Then the *rank* of q is defined to be rank(q) = P + N, and the *signature* of q is defined to be signature(q) = P - N.

Theorem 7.11 (Sylvester's Law of Inertia). Let $q(x) = x^T A x$ be a quadratic form, and suppose it has two diagonalised forms x = My and x = Nz, so $M^T A M = Diag(c_1, ..., c_n)$ and $N^T A N = Diag(d_1, ..., d_n)$. Then the ranks of $c_1y_1^2 + ... + c_ny_n^2$ and $d_1y_1^2 + ... + d_ny_n^2$ are equal, and their signatures are also equal.

Proof. This is omitted.

From the last theorem, it follows that we can define the rank of a quadratic form to be the rank of *any* quadratic form of it. Likewise for signature.

Example 7.12. In Example 7.8 we found a diagonal form $u^2 - 4v^2 + \frac{1}{4}w^2$. Thus, the rank is 2 + 1 = 3, and the signature is 2 - 1 = 1.

In Example 7.9, a diagonal form was $4v^2 - w^2$. The rank is 1 + 1 = 2, and the signature is 1 - 1 = 0.

Rank and signature have certain geometrical meaning.

Example 7.13. Consider the equation $cx^2 + dy^2 = 1$.

- If c, d > 0, this is an *ellipse*.
- If c > 0 and d < 0 then it is a *hyperbola*.
- If c > 0 and d = 0, it is the union of 2 lines (a *degenerate* case).

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