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1. Prelude

TODO: Fix notation here...

2. Introduction

In this course we builds up the rudiments of some important notions of algebraic structures. That is, a algebraic structure of an arbitrary set, or carrier set, coupled with various finitary operations defined on it. ..

3. Groups

Definition 3.1 (Binary operation). A binary operation on a set \mathcal{X} is a map $\circ : \mathcal{X} \times \mathcal{X} \to \mathcal{X}$. **N.B.** that the binary operation is *closed*.

Definition 3.2 (Magma). A **magma** is a set \mathcal{M} equipped with a binary operation \circ . We denote the magma as the tuple pair (\mathcal{M}, \circ) .

Definition 3.3 (Semi-group). A **semi-group** is a set \mathcal{G} equipped with binary operation that is associative. Hence, a semi-group is a magma where the operation is associative; That is, given any $x, y, z \in \mathcal{G}$ then $x \circ (y \circ z) = (x \circ y) \circ z \in \mathcal{G}$. We denote the semi-group as the tuple pair (\mathcal{G}, \circ) , not to be confused with a magma from context.

Definition 3.4 (Monoid). A semi-group with identity or, monoid for short, is a semi-group (\mathcal{G}, \circ) with a unique identity element $e \in \mathcal{G}$ such that $x \circ e = x = e \circ x \, \forall x \in \mathcal{G}$

Proof: unquieness of identity. Assume some other identity $e^{'}$ exists in \mathcal{G} then, $e^{'}=e^{'}\circ e=e\circ e^{'}=e$. \square

Example 3.5. Given $\mathcal{G} = \mathbb{N}$ with the binary law of composition \circ to be defined as arithmetic addition +. Then, $(\mathbb{N}, +)$ forms a semi-group with identity 0. Verify the axioms.

Definition 3.6 (Group). A **group** is a monoid where every element has an inverse. A abelian group is a group that is commutative.

Example 3.7. Given $\mathcal{G} = \mathbb{Z}$ with the binary law of composition \circ to be defined as arithmetic addition +. Then, $(\mathbb{Z}, +)$ forms a semi-group with identity 0. Verify the axioms.

Question 3.8. Why does the set of naturals \mathbb{N} not form a group under addition, however does form a monoid?

Definition 3.9 (Subgroup). A group \mathcal{H} is a **subgroup** of a group \mathcal{G} if the restriction of the binary operation $\circ: \mathcal{H} \times \mathcal{H} \to \mathcal{H}$ is a group operation on \mathcal{H} . In particular, A non-empty subset \mathcal{H} of a group \mathcal{G} is a subgroup of \mathcal{G} if and only if $h_1 \circ h_2 \in \mathcal{H}$ whenever $h_1, h_2 \in \mathcal{H}$, and $h^{-1} \in \mathcal{H}$ whenever $h \in \mathcal{H}$. We denote the subgroup by $\mathcal{H} \leq \mathcal{G}$.

Definition 3.10 (Smallest subgroup). If \mathcal{A} is a subset of a group \mathcal{G} , there is a *smallest* subgroup $Gp(\mathcal{A})$ of \mathcal{G} which contains \mathcal{A} , the subgroup *generated* by \mathcal{A} .

Example 3.11. Suppose $\mathcal{A} = \{g\}$ then $\operatorname{Gp}(\mathcal{A}) = \operatorname{Gp}(g)$ and so $\operatorname{Gp}(g) = \{g^n : n \in \mathbb{Z}\}$, where $g^0 = e, g^n$ is the product of n copies of g where n > 0, and g^n is the product of |n| copies of g^{-1} when n < 0.

Definition 3.12 (Cyclic group). A group \mathcal{G} is *cyclic* if $\mathcal{G} = \operatorname{Gp}(g)$ for some $g \in \mathcal{G}$. Such a element is called a *generator* of the group.

Definition 3.13 (Group order). If a group \mathcal{G} has finitely many elements, then the *order* $o(\mathcal{G})$ is the number of elements of \mathcal{G} .

Definition 3.14 (Normal subgroup). A subgroup \mathcal{H} of a group \mathcal{G} is a **normal**, or *self-conjugate*, if $g^{-1}hg \in \mathcal{H} \, \forall g \in \mathcal{G}$ and $h \in \mathcal{H}$. We denote the normal $\mathcal{H} \subseteq \mathcal{G}$.

Definition 3.15 (Simple group). A group \mathcal{G} is **simple** if it has no normal subgroups other than $\{e\}$ and \mathcal{G} .

3.1. **Group Homomorphisms.** Homomorphisms are structure preserving mappings. In group homomorphisms we preserve the structure of the binary operation ∘ as follows:

Definition 3.16 (Homomorphism). Let \mathcal{G} and \mathcal{H} be two groups. Then a mapping

$$\varphi:\mathcal{G}\to\mathcal{H}$$

is called a homomorphism if

$$\varphi(x \circ y) = \varphi(x) \circ \varphi(y) : x, y \in \mathcal{G}$$

It follows that, for some $g \in \mathcal{G}$ we have,

$$\varphi(e_g) = \varphi(g \circ g^{-1})$$

$$= \varphi(g) \circ \varphi(g^{-1})$$

$$= \varphi(g) \circ (\varphi(g))^{-1}$$

$$= e_h \in \mathcal{H}.$$

That is the identity e has been preserved. Hence, it does not matter if we compose in \mathcal{G} and map to \mathcal{H} or take two elements in \mathcal{G} then compose the mapped elements in \mathcal{H} , since the group structure has been preserved.

How much information about the elements inside the structure is, however, another quality to consider. Hence we fix some terminology here.

- A homomorphism that is injective is called monomorphic.
- A homomorphism that is surjective is called epimorphic.

• A homomorphism that is bijective is called isomorphic.

Thus we have the following definitions by considering a group homomorphism $\varphi: \mathcal{G} \to \mathcal{H}$.

Definition 3.17 (Monomorphic). φ is **monomorphic** if for $\varphi(x) = \varphi(y) \implies x = y \, \forall x, y \in \mathcal{G}$.

Definition 3.18 (Epimorphic). φ is **epimorphic** if $\forall h \in \mathcal{H} \exists g \in \mathcal{G}$ so that $\varphi(g) = h$.

Definition 3.19 (Isomorphic). φ is **isomorphic** if φ is **both** mono- and epic- morphic.

Some special cases are sometimes of particular interest and we shall outline them now.

Definition 3.20 (Endomorphic). A monomorphism $\mathcal{G} \to \mathcal{G}$ for a group \mathcal{G} is called an *endomorphism* of \mathcal{G} .

Definition 3.21 (Automorphic). A isomorphism $\mathcal{G} \to \mathcal{G}$ for a group \mathcal{G} is called an *automorphism* of \mathcal{G} .

Remark. The set $Aut(\mathcal{G})$ of automorphisms of \mathcal{G} forms a group, when composition of mappings is taken as the group law of composition.

3.2. Properties of homomorphisms.

Definition 3.22 (kernel). If $\varphi : \mathcal{G} \to \mathcal{H}$ is a group homomorphism, then the *kernel* is the set $\ker(\varphi) = \{g \in \mathcal{G} : \varphi(g) = e_h \in \mathcal{H}\}.$

If $\varphi: \mathcal{G} \to \mathcal{H}$ is a group homomorphism, then observe that $\ker(\varphi)$ is a normal subgroup of of \mathcal{G} .

3.3. Cosets. Let \mathcal{G} be a group and \mathcal{H} be a subgroup of \mathcal{G} with $g \in \mathcal{G} : g \notin \mathcal{H}$, then

Definition 3.23 (Left Coset). $gH = \{gh : h \in H\}$ is a **left coset of** \mathcal{H} in \mathcal{G} .

Definition 3.24 (Right Coset). $Hg = \{hg : h \in H\}$ is a **right coset of** \mathcal{H} in \mathcal{G} .

Definition 3.25 (Normal Subgroup). If gH = Hg then \mathcal{H} is a **normal** subgroup of \mathcal{G} , denoted by $\mathcal{H} \triangleleft \mathcal{G}$.

3.4. Factor (or Quotient) groups. Let \mathcal{G} be a commutative group and consider a subgroup \mathcal{H} . Then \mathcal{H} determines an equivalence relation in \mathcal{G} given by

$$x \sim x' \text{ iff } x - x' \in \mathcal{H}.$$

..

3.5. **Non-commutative Groups.** A common class of non-commutative groups are transformation groups. Note:

Definition 3.26 (Transformation). A bijective map $\varphi: X \to X$ is called a **transformation** of X. Note. The most trivial case is the *identity map* id_X by $id_X(x) = x$, $\forall x \in X$.

Hence, there exists a inverse φ^{-1} of φ such that $\varphi^{-1} \circ \varphi = id_X = \varphi \circ \varphi^{-1}$. Now, take two transformations of X, φ and ψ , and let the product $\varphi \circ \psi$ be well defined. Then the set of all transformations of X form the group **Transf(X)**. Since, given $\varphi, \psi, \phi \in Transf(X)$ then we have associativity, $\varphi \circ (\psi \circ \phi) = (\varphi \circ \psi) \circ \phi$. We have identity $e = id_X \in Transf(X)$ and so, inverses $\forall \varphi \in Transf(X) \exists ! \varphi^{-1} : \varphi \circ \varphi^{-1} = e$. Closure follows from the composition of two transformations φ and ψ , since $(\varphi \circ \psi)^{-1} = \psi^{-1} \circ \varphi^{-1}$.

A transformation group is a type of group action which describes symmetries of objects. More abstractly, since a group \mathcal{G} is a category with a single object in which every morphism is bijective. Then, a group action is a *forgetful functor* \mathcal{F} from the group \mathcal{G} in the category **Grp** to the set category **Set** that is, $\mathcal{F}: \mathcal{G} \to \mathbf{Set}$.

3.6. **Group actions.** For any mathematical object (e.g. sets, groups, vector spaces) X a isomorphism of X is a symmetry of X. The set of all isomorphisms of X, or symmetries of X, form a group called the symmetry group of X, denoted Sym(X). More formally;

Definition 3.27 (Group action). An action of a group \mathcal{G} on a mathematical object X is a mapping $\mathcal{G} \times X \to X$, defined by $(g, x) \mapsto g\dot{x}$ satisfying:

- $e.x = x \, \forall x \in X$ and
- $(gh).x = g.(g.x) \forall g, h \in \mathcal{G}, x \in X.$

That is, we have the (left) \mathcal{G} -action on X and denote this by $\mathcal{G} \curvearrowright X$.

Notice that we may study properties of the symmetries of some mathematical object X without reference to the structure of X in particular.

- 3.7. **Permutations.** Take a finite set X with |X| = n, then the transformations of X are called **permutations** of the elements of X. In particular, the group of permutations of $X = \{1, 2, \dots, n\}$ is a **symmetric group**, denoted S_n , with **order** $|S_n| = n!$. Thus, by taking any subgroup of S_n we have a **permutation group**. Also note that, for finite sets, *permutation* and *bijective maps* refer to the same operation, namely rearrangement of elements of X. Another way is to consider, a group \mathcal{G} and set X. Then a group action is defined as a group homomorphism φ from \mathcal{G} to the symmetric group of X. That is, the action $\varphi : \mathcal{G} \to S_n(X)$, assigns a permutation of X to each element of the group \mathcal{G} in the following way:
 - From the identity element $e \in \mathcal{G}$ to the identity transformation id_X of X, that is, $\varphi : e \to id_X$;

• A product of group homomorphisms $\varphi \circ \psi \in \mathcal{G}$ is then the composite of permutations given by φ and ψ in X.

Given that each element of \mathcal{G} is represented as a permutation. Then a group action can also be consider as a permutation representation.

A permutation $\sigma \in S_n$ can be written,

$$\sigma = \begin{pmatrix} 1 & 2 & \cdots & n \\ a_1 & a_2 & \cdots & a_n \end{pmatrix} \text{ where } a_1 = \sigma(1), a_2 = \sigma(2), \cdots.$$

The identity permutation $id_n \in S_n$ is simply,

$$id_n = \begin{pmatrix} 1 & 2 & \cdots & n \\ 1 & 2 & \cdots & n \end{pmatrix}$$

Since $|S_n| = n!$ then the total number of ways n elements maybe permuted is n!.

Take any two permutations $\sigma, \pi \in S_n$ then composition is well defined as **functional composition** as follows.

Given,

$$\sigma = \begin{pmatrix} 1 & 2 & \cdots & n \\ a_1 & a_2 & \cdots & a_n \end{pmatrix} \text{ and } \pi = \begin{pmatrix} a_1 & a_2 & \cdots & a_n \\ b_1 & b_2 & \cdots & b_n \end{pmatrix}$$

then,

$$\pi \circ \sigma = \begin{pmatrix} 1 & 2 & \cdots & n \\ \pi(a_1) & \pi(a_2) & \cdots & \pi(a_n) \end{pmatrix}$$
$$= \begin{pmatrix} 1 & 2 & \cdots & n \\ b_1 & b_2 & \cdots & b_n \end{pmatrix}$$

A inverse of any permutation $\sigma \in S_n$ is given by,

$$\sigma^{-1} = \begin{pmatrix} 1 & 2 & \cdots & n \\ a_1 & a_2 & \cdots & a_n \end{pmatrix}^{-1} = \begin{pmatrix} a_1 & a_2 & \cdots & a_n \\ 1 & 2 & \cdots & n \end{pmatrix}$$

3.8. **Permutation parity.** Consider the algebraic structure:

$$\triangle_n(x_1,\ldots,x_n) = \prod_{i< j} (x_i - x_j)$$

TODO..

3.9. Fields. We now may build higher order algebraic structures using the notion of a group.

Definition 3.28 (Field). A field \mathbb{F} is a set together with two binary operations, addition and multiplication, such that:

- addition forms a abelian group,
- multiplication forms a abelian quasi-group, i.e. a commutative multiplicative group on the set $\mathbb{F} \{0\}$,

coupled together with a law of distribution between the two binary operations.

4. Exact sequence

An **exact sequence** may either be a finite or infinite sequence of objects and morphisms between them. Such a sequence is constructed so that the image of one morphism equals the kernel of the next.

In particular;

Definition 4.1 (Exact Sequence). Consider the sequence of n group homomorphism between n+1 groups in the following way:

$$\mathcal{G}_0 \xrightarrow{\varphi_1} \mathcal{G}_1 \xrightarrow{\varphi_2} \mathcal{G}_2 \xrightarrow{\varphi_3} \dots \xrightarrow{\varphi_n} \mathcal{G}_n$$

Then the sequence is said to be exact if,

$$\ker(\varphi_{k+1}) = \operatorname{im}(\varphi_k)$$

for every $k \in \{1 \dots n\}$. For n = 3 the sequence is said to be a **short exact sequence**.

Example 4.2. Suppose we have $\mathcal{K} \subseteq \mathcal{G}$ and that $q: \mathcal{G} \to \mathcal{G}/\mathcal{K}$ is the quotient mapping. Then,

$$1 \longrightarrow \mathcal{K} \stackrel{\subseteq}{\longrightarrow} \mathcal{G} \stackrel{q}{\longrightarrow} \mathcal{G}/\mathcal{K} \longrightarrow 1$$

is a short exact sequence.

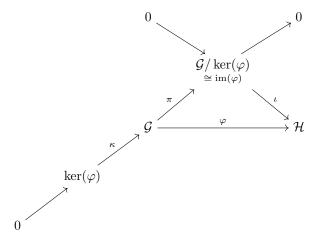
5. First Isomorphism Theorem

Definition 5.1. Let \mathcal{G} and \mathcal{H} , and let $\varphi: \mathcal{G} \to \mathcal{H}$ be a group homomorphism. Then:

- The kernel of φ is a normal subgroup of \mathcal{G} ; $\ker(\varphi) \leq \mathcal{G}$,
- The image of φ is a subgroup of \mathcal{H} ; $\operatorname{im}(\varphi) \leq \mathcal{H}$, and
- The image of φ is also isomorphic to the factor group $\mathcal{G}/\ker(\varphi)$; $\operatorname{im}(\varphi) \cong \mathcal{G}/\ker(\varphi)$.

In particular, if φ is epimorphic then $\mathcal{H} \cong \mathcal{G}/\ker(\varphi)$.

We may represent these fundamental relations in the following commutative diagram.



Notice the *exact sequence* that runs from the lower left to the upper right of the commutative diagram.

6. Determinants

Definition 6.1 (Determinant).

$$\det(A) = \sum_{\sigma \in S_n} sgn(\sigma) \prod_{i=1}^n a_{i,\sigma(i)}.$$

Example 6.2. Find the determinant of some matrix $A \in \mathcal{M}_{2,2}(\mathbb{F})$.

Suppose

$$A = \begin{pmatrix} a_{11} & a_{12} \\ a_{21} & a_{22} \end{pmatrix}$$

and take the permutation group

$$S_2 = \{\sigma_1, \sigma_2\}$$

with $\sigma_1 = id$ and $\sigma_2 = [1\,2]$. Then, by definition, we have:

$$\det(A) = \sum_{\sigma \in S_2} sgn(\sigma) \prod_{i=1}^2 a_{i,\sigma(i)}$$

$$= sgn(\sigma_1)(a_{1,\sigma_1(1)}a_{2,\sigma_1(2)}) + sgn(\sigma_2)(a_{1,\sigma_2(1)}a_{2,\sigma_2(2)})$$

$$= +(a_{11}a_{22}) - (a_{12}a_{21})$$

$$= a_{11}a_{22} - a_{12}a_{21}.$$

Example 6.3. Find the determinant of some matrix $A \in \mathcal{M}_{3,3}(\mathbb{F})$.

Suppose

$$A = \begin{pmatrix} a_{11} & a_{12} & a_{13} \\ a_{21} & a_{22} & a_{23} \\ a_{31} & a_{32} & a_{33} \end{pmatrix}$$

and take the permutation group

$$S_6 = \{\sigma_1, \sigma_2, \sigma_3, \sigma_4, \sigma_5, \sigma_6\}$$

with

$$\sigma_1 = [1 \ 2 \ 3] \ \sigma_2 = [1 \ 3 \ 2]$$

$$\sigma_3 = [2 \ 1 \ 3] \ \sigma_4 = [2 \ 3 \ 1]$$

$$\sigma_5 = [3 \ 2 \ 1] \ \sigma_6 = [3 \ 1 \ 2]$$

Then, by definition, we have:

$$\begin{split} \det(A) &= \sum_{\sigma \in S_6} sgn(\sigma) \prod_{i=1}^6 a_{i,\sigma(i)} \\ &= sgn(\sigma_1) (a_{1,\sigma_1(1)} a_{2,\sigma_1(2)} a_{3,\sigma_1(3)} a_{4,\sigma_1(4)} a_{5,\sigma_1(5)} a_{6,\sigma_1(6)}) \\ &+ sgn(\sigma_2) (a_{1,\sigma_2(1)} a_{2,\sigma_2(2)} a_{3,\sigma_2(3)} a_{4,\sigma_2(4)} a_{5,\sigma_2(5)} a_{6,\sigma_2(6)}) \\ &+ sgn(\sigma_3) (a_{1,\sigma_3(1)} a_{2,\sigma_3(2)} a_{3,\sigma_3(3)} a_{4,\sigma_3(4)} a_{5,\sigma_3(5)} a_{6,\sigma_3(6)}) \\ &+ sgn(\sigma_4) (a_{1,\sigma_4(1)} a_{2,\sigma_4(2)} a_{3,\sigma_4(3)} a_{4,\sigma_4(4)} a_{5,\sigma_4(5)} a_{6,\sigma_4(6)}) \\ &+ sgn(\sigma_5) (a_{1,\sigma_5(1)} a_{2,\sigma_5(2)} a_{3,\sigma_5(3)} a_{4,\sigma_5(4)} a_{5,\sigma_5(5)} a_{6,\sigma_5(6)}) \\ &+ sgn(\sigma_6) (a_{1,\sigma_6(1)} a_{2,\sigma_6(2)} a_{3,\sigma_6(3)} a_{4,\sigma_6(4)} a_{5,\sigma_6(5)} a_{6,\sigma_6(6)}) \end{split}$$

Lemma 6.4.

$$\det(A^T) = \det(A).$$

Proof. Consider some σ such that $\sigma_1, \ldots, \sigma_n$ is in order so that $\tau = \sigma^{-1}$ and take the matrix A_j^i so that $(A_j^i)^T = A_i^j$.

Then we see,

$$A_j^{\sigma(i)} = A_{\sigma^{-1}(j)}^j = A_{\tau(j)}^j$$

with $sgn(\tau) = sgn(\sigma)$. So,

$$\det(A) = \sum_{\sigma \in S_n} sgn(\sigma) \prod_{i=1}^n a_{i,\sigma(i)}$$

$$= \sum_{\tau \in S_n} sgn(\tau) \prod_{i=1}^n a_{\tau(i),\tau \cdot \sigma(i)}$$

$$= \sum_{\tau \in S_n} sgn(\tau) \prod_{i=1}^n a_{\tau(i),i}$$

$$= \det(A^T).$$

Lemma 6.5.

$$\det(AB) = \det(A)\det(B).$$

Proof. Let
$$A = [a_j^i]$$
 and $B = [b_j^i]$ with $A, B \in \mathcal{M}_{nn}(\mathbb{F})$ so that $AB = \left[\sum_{k=1}^n a_k^i b_j^k\right]$.

$$\det(AB) = \sum_{\sigma \in S_n} sgn(\sigma) \prod_{i=1}^n \left[\sum_{k=1}^n a_k^i b_{\sigma(i)}^k \right]$$

7. Adjugate (Classical Adjoint)

Definition 7.1 (Cofactor). For some square matrix $A \in \mathcal{M}_{n,n}(\mathbb{F})$ then the minor of the a_{ij} entry, denoted by $M_{i,j}$, is defined to be the determinant of the *submatrix* obtained by removing the $(i,j)^{th}$ row and column from A. That is,

$$C_{i,j} = (-1)^{i+j} M_{i,j}$$

where $C_{i,j}$ is called the *cofactor* of $a_{i,j}$.

Example 7.2. Consider the matrix $A \in \mathcal{M}_{3,3}$ and suppose

$$A = \begin{pmatrix} a_{11} & a_{12} & a_{13} \\ a_{21} & a_{22} & a_{23} \\ a_{31} & a_{32} & a_{33} \end{pmatrix}.$$

We may find the $C_{2,3}$ cofactor in the following way;

Observe that the minor $M_{2,3}$ may be found as follows,

$$M_{2,3} = \begin{vmatrix} a_{11} & a_{12} \\ a_{31} & a_{32} \end{vmatrix}$$
$$= a_{11}a_{32} - a_{12}a_{31}$$

and so the cofactor $C_{2,3}$ is by definition,

$$C_{2,3} = (-1)^{2+3} M_{2,3}$$

= $(-1)^5 (a_{11}a_{32} - a_{12}a_{31})$
= $a_{12}a_{31} - a_{11}a_{31}$.

Notice that we may now find the determinant of some $n \times n$ square matrix in terms of its cofactors. The process of this cofactor expansion is called the Laplace expansion.

Theorem 7.3. Suppose $A \in \mathcal{M}_{n,n}$ is some $n \times n$ square matrix then

in terms of the ith row we have,

$$\det(A) = \sum_{i=1}^{n} a_{ij} C_{ij}$$

or in terms of the j^{th} column

$$=\sum_{i=1}^{n}a_{ij}C_{ij}.$$

Proof. TODO: Prove Laplace expansion.

Example 7.4. Consider the matrix

$$A = \begin{pmatrix} 1 & 3 & 5 \\ 2 & 1 & 3 \\ 5 & 2 & 1 \end{pmatrix}.$$

Then,

$$\det(A) = 1 \cdot \begin{vmatrix} 1 & 3 \\ 2 & 1 \end{vmatrix} - 3 \cdot \begin{vmatrix} 2 & 3 \\ 5 & 1 \end{vmatrix} + 5 \cdot \begin{vmatrix} 2 & 1 \\ 5 & 2 \end{vmatrix}$$
$$= 1 \cdot (-5) - 3 \cdot (-13) + 5 \cdot (-1) = 29.$$

Recall that the determinant is a measure of overall scaling of a matrix. Observe that each cofactor are essentially sub determinants or sub scalings. By this we are motivated to form a matrix of these cofactors in the following way:

Definition 7.5 (Matrix of cofactors). Suppose $A \in \mathcal{M}_{n,n}(\mathbb{R})$ is some $n \times n$ real square matrix given by,

$$A = \begin{pmatrix} a_{11} & a_{12} & \cdots & a_{1n} \\ a_{21} & a_{22} & \cdots & a_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{n1} & a_{n2} & \cdots & a_{nn} \end{pmatrix}$$

with cofactors C_{ij} of a_{ij} . Then the cofactor matrix C is defined as,

$$C = \begin{pmatrix} C_{11} & C_{12} & \cdots & C_{1n} \\ C_{21} & C_{22} & \cdots & C_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ C_{n1} & C_{n2} & \cdots & C_{nn} \end{pmatrix}.$$

Now since a matrix of a linear transformation is really just a linear dilation (Invariant Subspaces) of each dimension of the coordinate system (Invariant Subspaces). Then for a real matrix we may find the inverse by means of dilating each respective dimension back. This is essentially what the adjugate matrix is. In fact, even when the matrix is *singular* (non-invertible) the adjugate is still well defined. That is, the adjugate is in actual fact, the pre-image of a real linear transformation.

More formally, we have the following definition;

Definition 7.6 (Adjugate (Classical adjoint)). Suppose $A \in \mathcal{M}_{n,n}(\mathbb{R})$ is some real $n \times n$ square matrix with cofactor matrix C. Then the *adjugate* or classical adjoint of A is defined as,

$$adj(A) = C^T$$
.

Thus we have the following result:

Lemma 7.7. For some non-singular $n \times n$ real square matrix $A \in \mathcal{M}_{n,n}(\mathbb{R})$ we have,

$$adj(A) = A^{-1} \det(A).$$

Proof. TODO: prove this lemma..

Example 7.8.

Suppose $A \in \mathcal{M}_{2,2}(\mathbb{R})$ is given by,

$$A = \begin{pmatrix} a_{11} & a_{12} \\ a_{21} & a_{22} \end{pmatrix}$$

then the adjugate of A is given by,

$$adj(A) = \begin{pmatrix} a_{22} & -a_{21} \\ -a_{12} & a_{11} \end{pmatrix}^{T}$$
$$= \begin{pmatrix} a_{22} & -a_{12} \\ -a_{21} & a_{11} \end{pmatrix}.$$

Some useful properties follow from the adjugate:

- adj(I) = I,
- adj(AB) = adj(B)adj(A),
- $adj(A^T) = adj(A)^T$.

Proof. TODO prove the above properties.

8. Vector Spaces

Definition 8.1 (Vector space). A vector space \mathcal{V} over a field \mathbb{F} is an abelian additive group $(\mathcal{V}, +)$, coupled with a mapping (scalar multiplication) $(\lambda, x) \to \lambda x$ of $\mathbb{F} \times \mathcal{V}$ into \mathcal{V} which satisfies

- 1.x = x.
- $(\lambda + \mu)(x) = \lambda x + \mu x$,
- $\lambda(\mu x) = (\lambda \mu)x$,
- $\lambda(x+y) = \lambda x + \lambda y$,

for $\lambda, \mu \in \mathbb{F}$ and $x, y \in \mathcal{V}$. We note here that the elements of \mathcal{V} are called *vectors* and that the elements of \mathbb{F} are called *scalars*.

Remark. Note that any field \mathbb{F} forms a degenerate vector space over itself.

Example 8.2. For a set X, the set of functions from $X \to \mathbb{F}$, denoted $Fun(X, \mathbb{F})$, is a vector space with pointwise addition and scalar multiplication. That is, for some $\varphi, \psi \in Fun(X, \mathbb{F})$, $x \in X$ and $\lambda \in \mathbb{F}$ we have,

$$(\varphi + \psi)(x) = \varphi(x) + \psi(x)$$

and

$$(\lambda \varphi)(x) = \lambda(\varphi(x)).$$

8.1. Subspaces.

Definition 8.3. A non-empty subset \mathcal{U} of a vector space \mathcal{V} over \mathbb{F} is a *linear subspace* if it is a subgroup of \mathcal{V} and if $\lambda x \in \mathcal{U}$ whenever $\lambda \in \mathbb{F}$ and $x \in \mathcal{U}$ written $\mathcal{U} \leq \mathcal{V}$.

Corollary. A non-empty subset \mathcal{U} of a vector space \mathcal{V} over a field \mathbb{F} is a linear subspace if and only if, for any $u, u' \in \mathcal{U}$ and $\lambda \in \mathbb{F}$ we have, $u + \lambda u' \in \mathcal{U}$.

Example 8.4. Let the set of all real-valued functions on \mathbb{R}^n be denoted by $\mathcal{F}(\mathbb{R}^n)$. Notice that, for any $f, g \in \mathcal{F}(\mathbb{R}^n)$, $\mathbf{x} \in \mathbb{R}^n$ and $\lambda \in \mathbb{R}$ we have, $(f+g)(\mathbf{x}) = f(\mathbf{x}) + g(\mathbf{x})$ and $(\lambda f)(\mathbf{x}) = \lambda f(\mathbf{x})$. That is we claim that, $\mathcal{F}(\mathbb{R}^n)$ is a vector space.

Now, observe that the set of all continuous real-valued functions, denoted $\mathcal{C}(\mathbb{R}^n)$ is a subspace of $\mathcal{F}(\mathbb{R}^n)$.

Example 8.5. Given two subspaces $\mathcal{U}, \mathcal{V} \leq \mathcal{W}$ of vector space \mathcal{W} , then $\mathcal{U} \cap \mathcal{V} \leq \mathcal{W}$.

Proof.

Suppose $x, \lambda x' \in \mathcal{U} \cap \mathcal{V}$ with any $\lambda \in \mathbb{F}$. Then,

$$x, \lambda x' \in \mathcal{U}$$
 and $x, \lambda x' \in \mathcal{V}$.

Now, since both \mathcal{U} and \mathcal{V} are subspaces then,

$$x + \lambda x' \in \mathcal{U} \text{ and } x + \lambda x' \in \mathcal{V}$$

 $\implies x + \lambda x' \in \mathcal{U} \cap \mathcal{V}.$

Lemma 8.6. In general, for a vector space V and any family of subspaces $U_i \leq V$ for every $i \in I$, we have $\bigcap_{i \in I} U_i \leq V$.

Example 8.7. Consider the subspaces $\mathcal{V}', \mathcal{V}'', \mathcal{W} \leq \mathcal{V}$. Prove that

$$(\mathcal{V}' \cap \mathcal{W}) + (\mathcal{V}'' \cap \mathcal{W}) \subseteq (\mathcal{V}' + \mathcal{V}'') \cap \mathcal{W}.$$

Proof.

Pick some

$$x \in (\mathcal{V}' \cap \mathcal{W}) + (\mathcal{V}'' \cap \mathcal{W}) : x = u + v$$

where $u \in \mathcal{V}' \cap \mathcal{W}$ and $v \in \mathcal{V}'' \cap \mathcal{W}$.

Then,

$$u \in \mathcal{V}', u \in \mathcal{W}$$

and

$$v \in \mathcal{V}'', v \in \mathcal{W}.$$

Now, since $u, v \in \mathcal{W}$ and that \mathcal{W} is a subspace,

$$x = u + v \in \mathcal{W}$$
.

Also, since

$$u \in \mathcal{V}'$$
 and $v \in \mathcal{V}''$

then

$$u+v\in\mathcal{V}'+\mathcal{V}''$$
.

Hence,

$$x \in W \cap (\mathcal{V}' + \mathcal{V}'').$$

Does the reverse inclusion always hold?

Proof.

Pick some $y \in (\mathcal{V}' + \mathcal{V}'') \cap W$ so that

$$y \in (\mathcal{V}' + \mathcal{V}'')$$
 or $y \in \mathcal{W}$.

Now, write y = y' + y'' with $y' \in \mathcal{V}'$ and $y'' \in \mathcal{V}''$. Observe that,

for
$$y' \in \mathcal{W}$$
 so that $y' \in \mathcal{V}' \cap \mathcal{W}$
 $\Rightarrow y'' \in \mathcal{W}$ so that $y'' \in \mathcal{V}'' \cap \mathcal{W}$

and so,

$$y = y' + y'' \notin (\mathcal{V}' \cap \mathcal{W}) + (\mathcal{V}'' \cap \mathcal{W}).$$

Hence no. \Box

Example 8.8. Consider the subspaces $\mathcal{V}', \mathcal{V}'' \leq \mathcal{V}$ and suppose that $\mathcal{V}' \cup \mathcal{V}''$ is also a subspace of \mathcal{V} . Then show that either $\mathcal{V}' \subseteq \mathcal{V}''$ or $\mathcal{V}'' \subseteq \mathcal{V}'$.

Proof.

Let $x \in \mathcal{V}'$ and $y \in \mathcal{V}''$ so that $x \in \mathcal{V}' \cup \mathcal{V}''$ and $y \in \mathcal{V}' \cup \mathcal{V}''$.

Now, since $\mathcal{V}' \cup \mathcal{V}''$ is a subspace then,

$$x + y \in \mathcal{V}' \cup \mathcal{V}''$$

so,

$$x + y \in \mathcal{V}'$$
 or $x + y \in \mathcal{V}''$.

Case i.) If $x + y \in \mathcal{V}'$ while $x \in \mathcal{V}'$ and that, since \mathcal{V}' is a subspace then, $y \in \mathcal{V}'$. So,

$$y \in \mathcal{V}'' \implies y \in \mathcal{V}'$$
$$\mathcal{V}' \subseteq \mathcal{V}''.$$

Case ii.) If $x + y \in \mathcal{V}''$ while $y \in \mathcal{V}''$ and that, since \mathcal{V}'' is a subspace then, $x \in \mathcal{V}''$. So,

$$x \in \mathcal{V}' \implies x \in \mathcal{V}''$$

$$\mathcal{V}'' \subseteq \mathcal{V}'.$$

Remark. Observe that, in general, unions of subspaces are not necessarily subspaces.

Example 8.9. Consider the subspaces $W, W' \leq V$ of some vector space V. Show that, if W'' is any subspace of V containing W and W' then $W + W' \subseteq W''$.

Proof.

If,

$$\mathcal{W} \cup \mathcal{W}' \subseteq \mathcal{W}''$$
 with $\mathcal{W}'' \leq \mathcal{V}$

then,

$$\mathcal{W} \subseteq \mathcal{W}''$$
 and $\mathcal{W}' \subseteq \mathcal{W}''$.

Now, suppose x = x' + x'' with $x' \in \mathcal{W}$ and $x'' \in \mathcal{W}'$. Then,

$$x \in \mathcal{W} \cup \mathcal{W}' \implies x \in \mathcal{W}''.$$

Hence,

$$W + W' \subseteq W''$$
.

9. Direct Sum

Consider two subspaces $\mathcal{U}, \mathcal{V} \leq \mathcal{W}$ of \mathcal{W} . Then \mathcal{U} and \mathcal{V} are said to be **complementry** if every vector $w \in \mathcal{W}$ has a *unique* decomposition $w = u + v : u \in \mathcal{U}$ and $v \in \mathcal{V}$. The vector space \mathcal{W} is then said to be the **internal direct sum** of subspaces \mathcal{U} and \mathcal{V} , written $\mathcal{W} = \mathcal{U} \oplus \mathcal{V}$.

Definition 9.1 (Direct sum). Consider two vector spaces \mathcal{U} and \mathcal{V} . Then the sum $\mathcal{U} + \mathcal{V}$ is called direct if and only if $\mathcal{U} \cap \mathcal{V} = \{0\}$, written $\mathcal{U} \oplus \mathcal{V}$.

Definition 9.2 (Complementary subspaces). Both \mathcal{U} and \mathcal{V} are complementary subspaces of \mathcal{W} if and only if,

- $W = \mathcal{U} + \mathcal{V}$ and,
- $\mathcal{U} \cap \mathcal{V} = \{\mathbf{0}\}.$

That is, if $\mathcal{U}, \mathcal{V} \leq \mathcal{W}$ we have that $\mathcal{W} = \mathcal{U} \oplus \mathcal{V}$.

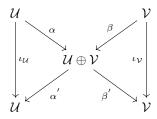
Definition 9.3 (External direct sum). Consider two arbitrary vector spaces \mathcal{U} and \mathcal{V} over some field \mathbb{F} . We may define their direct sum constructively by setting $\mathcal{U} \oplus \mathcal{V} = \mathcal{U} \times \mathcal{V} = \{(u, v) : u \in \mathcal{U}, v \in \mathcal{V}\}$ with vector addition and scalar multiplication defined by:

- (u, v) + (u', v') = (u + u', v + v') and
- $\lambda(u,v) = (\lambda u, \lambda v) : \lambda \in \mathbb{F}$.

Consider two vector spaces \mathcal{U} and \mathcal{V} with the linear map $\varphi:\mathcal{U}\to\mathcal{V}$ and the short exact sequence

$$\mathcal{U} \stackrel{\alpha}{\to} \mathcal{U}/\ker(\varphi) \stackrel{\gamma}{\to} \operatorname{im}(\varphi) \stackrel{\beta}{\to} \mathcal{V}.$$

In particular, observe that the map $\alpha: u \mapsto (u,0)$ and $\beta: v \mapsto (0,v)$ are injective. While the maps $\alpha': (u,0) \mapsto u$ and $\beta': (0,v) \mapsto v$ are subjective. That is γ is isomorphic. Hence we have the following commutative diagram:



Observe that the internal and external direct sums are isomorphic in this case.

We can now write higher dimensional vector spaces are direct sums of one dimensional subspaces by the natural isomorphism outlined above.

Example 9.4. Let $\mathbb{R}_1 = \{(\mathbf{x}_1, \mathbf{0}) : \mathbf{x}_1 \in \mathbb{R}\}$ and $\mathbb{R}_2 = \{(\mathbf{0}, \mathbf{x}_2) : \mathbf{x}_2 \in \mathbb{R}\}$. Then we see that $\mathbb{R}^2 = \mathbb{R}_1 \oplus \mathbb{R}_2 \cong \mathbb{R} \oplus \mathbb{R}$.

More generally, consider the subspace \mathbb{R}_i with vectors of the form

$$\mathbb{R}_i = \{(\mathbf{0}_1, \dots, \mathbf{0}_{i-1}, \mathbf{x}_i, \mathbf{0}_{i+1}, \dots, \mathbf{0}_n) : \mathbf{x}_i \in \mathbb{R}\}.$$

Then we have the following:

$$\mathbb{R}^{n} = \bigoplus_{i=1}^{n} \mathbb{R}_{i}$$

$$= \mathbb{R}_{1} \oplus \mathbb{R}_{2} \oplus \cdots \oplus \mathbb{R}_{n}$$

$$\cong \underbrace{\mathbb{R} \oplus \cdots \oplus \mathbb{R}}_{n \text{ times}}$$

$$= \bigoplus_{i=1}^{n} \mathbb{R}.$$

A subspace of \mathbb{R}^n that is isomorphic to \mathbb{R}_i is called the *canonical copy* of \mathbb{R}_i in \mathbb{R} .

Example 9.5.

Let,

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

$$B = \begin{pmatrix} 2 & 1 & 0 \\ 0 & 2 & 1 \\ 0 & 0 & 2 \end{pmatrix} \text{ and } C$$
 = (1).

Then,

$$A \oplus B \oplus C = \begin{pmatrix} A & 0 & 0 \\ 0 & B & 0 \\ 0 & 0 & C \end{pmatrix}.$$

Remark. Note that direct sums are associative, that is, $(A \oplus B) \oplus C = A \oplus (B \oplus C)$.

10. Invariant Subspaces

Definition 10.1 (Characteristic polynomial). Consider some matrix $A \in \mathcal{M}_{n,n}(\mathbb{F})$ with the eigenvector \mathbf{x} and eigenvalue $\lambda \in \mathbb{F}$ relation,

$$A\mathbf{x} = \lambda\mathbf{x} : \mathbf{x} \neq \mathbf{0}$$

then we have,

$$A\mathbf{x} - \lambda \mathbf{x} = \mathbf{0}$$
$$\Rightarrow (A - \lambda I_n)\mathbf{x} = \mathbf{0}.$$

Hence, the characteristic polynomial is given by

$$\chi_A(\lambda) = |A - \lambda I_n|.$$

Theorem 10.2 (Clayley-Hamilton). For any matrix $A \in \mathcal{M}_{n,n}(\mathbb{F})$ with characteristic polynomial $\chi_A(\lambda) : \lambda \in \mathbb{F}$. Then A satisfies $\chi_A(\lambda)$, that is, $\chi_A(A) = 0$.

Example 10.3. Suppose some matrix $A \in \mathcal{M}_{n,n}(\mathbb{F})$ is *idempotent*, that is $A^2 = A$. Then show that each eigenvalue of A is either zero or one.

Proof.

Suppose that,

$$A\mathbf{x} = \lambda \mathbf{x} : \mathbf{x} \neq \mathbf{0}$$

by induction we see that,

$$A^{2}\mathbf{x} = \lambda^{2}\mathbf{x}$$

$$\Rightarrow A\mathbf{x} = \lambda^{2}\mathbf{x}$$

$$= \lambda\mathbf{x}$$

$$\Rightarrow \lambda = \lambda^{2} \text{ iff } \lambda \in \{0, 1\}.$$

Example 10.4. Suppose some matrix $A \in \mathcal{M}_{n,n}(\mathbb{F})$ is *nilpotent*, that is $A^k = \mathbf{0}$ for some positive integer k. Then show that all the eigenvalues of A are zero.

Proof.

Suppose that,

$$A\mathbf{x} = \lambda\mathbf{x} : \mathbf{x} \neq \mathbf{0}$$

by induction we see that,

$$A^{k}\mathbf{x} = \lambda^{k}\mathbf{x}$$

$$\Rightarrow \mathbf{0} = \lambda^{k}\mathbf{x}$$

$$\Rightarrow \lambda^{k} = 0$$
and so, $\lambda = 0$ for every $\lambda \in \mathbb{F}$.

Definition 10.5 (Eigen Subspace). For some matrix $A \in \mathcal{M}_{n,n}$. Define a eigen-subspace \mathcal{E}_{λ} for every eigenvalue λ as the space that contains all the eigenvectors for that particular eigenvalue. That is,

$$\mathcal{E}_{\lambda}(k) = \ker(A - \lambda I_n)^k.$$

Definition 10.6 (Spectrum). The *spectrum* of a linear operator $A \in \mathcal{M}_{n,n}(\mathbb{C})$ is the set of all its eigenvalues, denoted by $\sigma(A)$. That is,

$$\sigma(A) = \{ \lambda \in \mathbb{C} : A\mathbf{x} = \lambda\mathbf{x} : \mathbf{x} \neq \mathbf{0} \}.$$

Definition 10.7 (Spectral radius). The spectral radius pf a linear operator $A \in \mathcal{M}_{n,n}(\mathbb{C})$ is the radius of the smallest disc centered at the origin in the complex plane that includes all the eigenvalues of A, denoted by $\rho(A)$. That is,

$$\rho(A) = \max\{|\lambda| : \lambda \in \sigma(A)\}.$$

Definition 10.8 (T-invariant subspace). For a vector space \mathcal{V} consider the linear endomorphism $T: \mathcal{V} \to \mathcal{V}$. Then, for some subspace $\mathcal{W} \leq \mathcal{V}$, \mathcal{W} is said to be T-invariant if $T(\mathcal{W}) \subseteq \mathcal{W}$.

Example 10.9. Show that,

$$W = span \left\{ \begin{pmatrix} 1\\1\\0\\1 \end{pmatrix}, \begin{pmatrix} 1\\0\\2\\0 \end{pmatrix} \right\} \text{ and } W' = span \left\{ \begin{pmatrix} 1\\1\\0\\0 \end{pmatrix}, \begin{pmatrix} 1\\1\\-1\\0 \end{pmatrix} \right\}$$

are T-invariant subspaces of the linear map:

$$T = \begin{pmatrix} 4 & -2 & -1 & -1 \\ 3 & -1 & -1 & -1 \\ -2 & 2 & 2 & 0 \\ 1 & -1 & 0 & 1 \end{pmatrix}.$$

Observe that,

$$T(W) = span \left\{ \begin{pmatrix} 1\\1\\0\\1 \end{pmatrix}, \begin{pmatrix} 2\\1\\2\\1 \end{pmatrix} \right\} \subseteq W$$

and that

$$T(W') = span \left\{ \begin{pmatrix} 2\\2\\0\\0 \end{pmatrix}, \begin{pmatrix} 3\\3\\-2\\0 \end{pmatrix} \right\} \subseteq W'$$

and so, both W and W' are T-invariant subspaces.

11. JORDAN CANONICAL FORM

The Jordan (normal) canonical form is a set of "almost diagonal" matries, called the *Jordan blocks*. The Jordan blocks contain the diagonal components composed of eigenvalues. We call this form, the *block diagonal form*. The set of Jordan blocks form a *equivalence class* under similarity invariance of the eigen-invariant subspaces.

Definition 11.1 (Jordan block). A Jordan block $\mathcal{J}_k(\lambda)$ is a $k \times k$ upper triangular square matrix of the form:

$$\mathcal{J}_k(\lambda) = \begin{pmatrix} \lambda & 1 & & & \\ 0 & \lambda & 1 & & 0 \\ & & \ddots & \ddots & \\ 0 & & & \ddots & 1 \\ & & & & \lambda \end{pmatrix}$$

where there are k-1 (+1)'s in the *superdiagonal* and k of a particular scalar eigenvalue λ in the main diagonal while all other entries are zero. The trivial 1×1 Jordan block of λ is $\mathcal{J}_1(\lambda) = (\lambda)$.

Definition 11.2 (Jordan matrix). A Jordan matrix $\mathcal{J} \in \mathcal{M}_{n,n}(\mathbb{F})$ is a direct sum (9) of Jordan blocks. That is,

$$\mathcal{J} = \bigoplus_{i=1}^{N} \mathcal{J}_{\lambda_i, m_k} : 1 \le i \le N.$$

where each $\lambda_i \in \sigma(A)$ for some matrix $A \in \mathcal{M}(\mathbb{F})$ need not be distinct.

Example 11.3.

Suppose,

$$A = \begin{pmatrix} 3 & 1 & 1 \\ -4 & -1 & 0 \\ 0 & 0 & 1 \end{pmatrix}.$$

We wish to find the Jordan canonical form so that we have the usual change of basis given by, $C\mathcal{J} = AC$. First we find the characteristic polynomial,

$$\chi_A(\lambda) = (3 - \lambda)(\lambda + 1)^2 - 4(\lambda + 1)$$

$$= -(\lambda + 1) \{(\lambda + 1)(\lambda - 3) + 4\}$$

$$= -(\lambda + 1)(\lambda^2 - 2\lambda + 1)$$

$$\Rightarrow \chi_A(\lambda) = (\lambda + 1)(\lambda - 1)^2$$

and so, $\lambda_1=-1$ and $\lambda_2=1$ with algebraic multiplicity two. Next recall that,

$$\mathcal{E}_{\lambda}(n) = \ker(A - \lambda I)^n$$

and, in particular, that

geometric multiplicity = dim $\mathcal{E}_{\lambda}(1)$.

For λ_1 we recall that the geometric multiplicity is always less than or equal to the algebraic multiplicity and since λ_1 has algebraic multiplicity of one it has geometric multiplicity of one, hence has the trivial Jordan block:

$$\mathcal{J}_1(\lambda_1) = \mathcal{J}_1(-1).$$

To find the corresponding eigenvector to λ_1 we look in the span of the kernel as follows

$$\mathbf{v}_{1} \in \mathcal{E}_{\lambda_{1}}(1) = \mathcal{E}_{-1}(1)$$

$$= \ker \left(A + I \right)^{1}$$

$$= \ker \left(\begin{pmatrix} 4 & 1 & 1 \\ -4 & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix} \right)$$

$$= \ker \left(\begin{pmatrix} 0 & 1 & 1 \\ 1 & 0 & 0 \\ 0 & 0 & 0 \end{pmatrix} \right)$$

$$\Rightarrow \mathbf{v}_{1} \in \operatorname{span} \left\{ \begin{pmatrix} 0 \\ 1 \\ -1 \end{pmatrix} \right\}.$$
(after a row reduction)

Note that it is important we find \mathbf{v}_1 first as we shall see. Now, see that

$$\mathcal{E}_{\lambda_{2}}(1) = \mathcal{E}_{1}(1) = \ker(A - I)^{1}$$

$$= \ker\begin{pmatrix} 2 & 1 & 1 \\ -4 & -2 & 0 \\ 0 & 0 & -2 \end{pmatrix}$$

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

$$= \ker\begin{pmatrix} 0 & 0 & 0 \\ 2 & 1 & 0 \\ 0 & 0 & 1 \end{pmatrix}$$
(after row reductions.)

Hence, dim $\mathcal{E}_{\lambda_2}(1) = 2$ and so we have the corresponding Jordan block,

$$\mathcal{J}_2(\lambda_2) = \mathcal{J}_2(1).$$

Now, to find the corresponding eigenvectors to λ_2 we need only pick a vector **not** in the span of $\mathcal{E}_{\lambda_2}(1)$ and is linearly independent from \mathbf{v}_1 . Consider the vector $\mathbf{v}_2 \notin \mathcal{E}_{\lambda_2}(1)$ with,

$$\mathbf{v}_2 = \begin{pmatrix} 0 \\ 1 \\ 0 \end{pmatrix}.$$

Then by the Jordan chain we see that,

$$\mathbf{v}_3 = \mathcal{E}_{\lambda_2}(1)\mathbf{v}_2$$

$$= \begin{pmatrix} 1 \\ -2 \\ 0 \end{pmatrix} \in span\mathcal{E}_{\lambda_2}(1)$$

thus, $\mathcal{E}_{\lambda_2}(1)\mathbf{v}_3 = \mathbf{0}$ and since the geometric multiplicity of λ_2 was two we now have the spanning set of corresponding eigenvectors $\mathbf{v}_2, \mathbf{v}_3$. So,

$$C\mathcal{J} = AC$$

$$\Rightarrow A = C\mathcal{J}C^{-1}$$

$$\Rightarrow A = C\mathcal{J}_{1}(\lambda_{1}) \oplus \mathcal{J}_{2}(\lambda_{2})C^{-1}$$

$$A = \begin{pmatrix} 0 & 0 & 1\\ 1 & 1 & -2\\ -1 & 0 & 0 \end{pmatrix} \begin{pmatrix} -1 & 0 & 0\\ 0 & 1 & 1\\ 0 & 0 & 1 \end{pmatrix} \begin{pmatrix} 0 & 0 & 1\\ 1 & 1 & -2\\ -1 & 0 & 0 \end{pmatrix}^{-1}.$$

12. Algebras

Definition 12.1 (Algebras). A finite-dimensional (associative) algebra \mathcal{A} over a field \mathbb{F} is a finite-dimensional vector space over \mathbb{F} equipped with a law of composition, that is, a mapping (multiplication) $(a,b) \to ab$ from $\mathcal{A} \times \mathcal{A}$ into \mathcal{A} which satisfies

- (ab)c = a(bc) (associativity),
- a(b+c) = ab + ac,
- (a+b)c = ac + bc,
- $\lambda(ab) = (\lambda a)b = a(\lambda b)$,

for $\lambda \in \mathbb{F}$ and $a, b, c \in \mathcal{A}$.

Remark. A algebra \mathcal{A} is called *unital* if there exists $1 \in \mathcal{A}$, the *identity element*, such that $1a = a1 = a \forall a \in \mathcal{A}$.

Remark. A algebra \mathcal{A} is commutative if $ab = ba \forall a, b \in \mathcal{A}$.

13. Vector space Homomorphisms

Definition 13.1 (Vector space Homomorphism). A vector space homomorphism is the structure preserving mapping between two vector spaces. That is, for some vector spaces \mathcal{U} and \mathcal{V} over common field \mathbb{F} with morphism $\varphi: \mathcal{U} \to \mathcal{V}$ we have,

$$\varphi(u_1 + \lambda u_2) = \varphi(u_1) + \lambda \varphi(u_2) : u_1, u_2 \in \mathcal{U} \text{ and } \lambda \in \mathbb{F}.$$

Remark. A vector space homomorphism preserves linearity and so we typically call the homomorphism a linear map.

Definition 13.2 (Kernel). Consider the linear map $\varphi : \mathcal{U} \to \mathcal{V}$, the *kernel* of φ , denoted $\ker(\varphi)$ is given by:

$$\ker(\varphi) = \{ u \in \mathcal{U} : \varphi(u) = \mathbf{0} \in \mathcal{V} \}.$$

The set or space of elements inside the kernel is thus called the *null space*.

Definition 13.3 (Nullity). The *nullality* is the dimension of the null space or kernel.

Remark. Note that the nullity gives a measure of the injectivity of the linear map φ .

Definition 13.4 (Image). For some morphism $\varphi: \mathcal{U} \to \mathcal{V}$ the *image* is defined as,

$$img(\varphi) = \varphi(\mathcal{U}) = \{v \in \mathcal{V} : v = \varphi(u) \text{ for some } u \in \mathcal{U}\}.$$

Definition 13.5 (Preimage). For some morphism $\varphi: \mathcal{U} \to \mathcal{V}$ the *preimage* is defined as,

$$\varphi$$
 $(\mathcal{V}) = \{ u \in \mathcal{U} : \varphi(u) \in \mathcal{V} \}$

Remark. The preimage can be defined even when no inverse morphism exists.

Theorem 13.6. Let \mathcal{U} and \mathcal{V} be vector spaces over some common field \mathbb{F} and some linear map $\varphi: \mathcal{U} \to \mathcal{V}$ then

$$\dim(img(\varphi)) + \dim(\ker(\varphi)) = \dim(\mathcal{U})$$

Definition 13.7 (Injective). If, for the linear map $\varphi : \mathcal{U} \to \mathcal{V}$, we have

$$\varphi(u_1) = \varphi(u_2) \implies u_1 = u_2 \,\forall u_1, u_2 \in \mathcal{U}$$

then φ is said to be *injective*.

Definition 13.8 (Surjective). If, for the linear map $\varphi : \mathcal{U} \to \mathcal{V}$, we have

$$\forall v \in \mathcal{V} \exists u \in \mathcal{U} : \varphi(u) = v$$

then φ is said to be *surjective*.

Definition 13.9 (Bijective). If the linear map $\varphi : \mathcal{U} \to \mathcal{V}$ is both injective and subjective then φ is said to be *bijective*, or, isomorphic in the context vector space homomorphisms.

13.1. **Properties of linear maps.** One may note that a linear map can always be represented in the form of a matrix and as such the determinant can be taken given that linear maps are from vector spaces to vector spaces over a common field. Hence, the conceptual question arises, geometrically, what exactly is the determinant? It turns out the determinant is a *measure* of dilation much like how scaler multiplies of the unit scalar dilate it. That is, 3*1=3 is the dilation of the unit scalar 1

by a measure of 3. Now, recall that a field is a vector space over itself and so we may use this to generalise our intuition from above. Thus, we may view a linear map as a linear dilation of some vector in one vector space to another vector space and the determinant as the measure of dilation of this mapping.

TODO.. motivate to prove the determinant definition here.. and motiave to the charteristic polynomial and eigen spaces..

14. Dual Spaces

Recall that; If \mathcal{U} and \mathcal{V} are two vector spaces over a field \mathbb{F} , then the set of all linear maps from $\mathcal{U} \to \mathcal{V}$, denoted $L(\mathcal{U}, \mathcal{V})$, is a vector space over \mathbb{F} .

Definition 14.1. Let \mathcal{V} be a vector space over \mathbb{F} . The dual space of \mathcal{V} , denoted \mathcal{V}^* , is the set of all linear maps from \mathcal{V} to \mathbb{F} . That is, $\mathcal{V}^* = L(\mathcal{V}, \mathbb{F})$. The elements of \mathcal{V}^* are called *linear functionals* on \mathcal{V} . Frequently, linear functionals on \mathcal{V} are called *covectors* or 1-forms.

Remark. Since $L(\mathcal{V}, \mathbb{F})$ is a vector space over \mathbb{F} then so too is the dual \mathcal{V}^* , with point-wise addition and scalar multiplication. Moreover, if $n = \dim \mathcal{V} < \infty$, then $\dim \mathcal{V}^* = \dim L(\mathcal{V}, \mathbb{F}) = n$.

Example 14.2. Fix $x \in [0,1]$ and define the evaluation map $ev_x : \mathcal{C}([0,1]) \to \mathbb{R}$ by $ev_x(f) = f(x)$ for any $f \in \mathcal{C}([0,1])$. Then $ev_x \in (\mathcal{C}([0,1]))^*$.

15. Normed vector space

A norm on a complex vector space \mathcal{V} is a real-valued function $\|\cdot\|: \mathcal{V} \to \mathbb{R}$ that, to any vector $x \in \mathcal{V}$ associates a real number $\|x\|$, satisfies the following:

- $||x|| \ge 0, \forall x \in \mathcal{V}$ and in particular $||x|| = 0 \Leftrightarrow x = 0$ (positive definite),
- $\|\alpha x\| = |\alpha| \|x\| \, \forall x \in \mathcal{V}$, and $\alpha \in \mathbb{C}$ (uniform scaling),
- $||x + y|| \le ||x|| + ||y|| \, \forall x, y \in \mathcal{V}$ (triangle inequality).

15.1. **Linear Operators.** In functional analysis one considers linear maps between normed vector spaces. In this context such linear mappings are termed *linear operators* or just operators. In this light one is motivated to talk about notions of *continuity* of such operators.

For any normed vector space $(\mathcal{U}, \|\cdot\|_{\mathcal{U}})$ consider the closed unit ball in \mathcal{U} given by $B(\mathcal{U}) = \{x \in \mathcal{U} : \|x\|_{\mathcal{U}} \leq 1\}$. Hence, we may now define boundedness of an operator formally as follows:

Definition 15.1 (Bounded linear operator). A linear operator $T: \mathcal{U} \to \mathcal{V}$ between normed vector spaces $(\mathcal{U}, \|\cdot\|_{\mathcal{U}})$ and $(\mathcal{V}, \|\cdot\|_{\mathcal{V}})$ is called a *bounded* linear operator if it maps the closed unit ball $B(\mathcal{U})$ in \mathcal{U} into a bounded set (a ball) in \mathcal{V} .

Corollary. A linear operator $T: \mathcal{U} \to \mathcal{V}$ between normed vector spaces $(\mathcal{U}, \|\cdot\|_{\mathcal{U}})$ and $(\mathcal{V}, \|\cdot\|_{\mathcal{V}})$ is continuous if and only if it is bounded.

16. Inner Product Spaces

Let V be a complex vector space.

Definition 16.1 (Inner product). An inner product in V is a mapping $\langle \cdot, \cdot \rangle : V \times V \to \mathbb{C}$ satisfying

- Hermitian symmetry: $\langle x, y \rangle = \overline{\langle y, x \rangle} \, \forall x, y \in V$;
- Linearity: $\langle \alpha x_1 + \beta x_2, y \rangle = \alpha \langle x_1, y \rangle + \beta \langle x_2, y \rangle \ \forall x_1, x_2, y \in V \ \text{and} \ \alpha, \beta \in \mathbb{C};$
- Positive definite: $\langle x, x \rangle \ge 0$ and $\langle x, x \rangle = 0$ iff $x = 0, \forall x \in V$;

Remark. Note that condition (i) implies that $\langle x, x \rangle \in \mathbb{R}$ and so condition (ii) is consistent.

Remark. Note that condition (ii) defines that the inner product is linear in the first argument. However, the inner product is *conjugate linear* in the second argument as shown by:

Proof.

$$\langle x, \alpha y_1 + \beta y_2 \rangle = \overline{\langle \alpha y_1 + \beta y_2, x \rangle}$$

$$= \overline{\langle \alpha y_1, x \rangle + \langle \beta y_2, x \rangle}$$

$$= \overline{\alpha} \overline{\langle y_1, x \rangle} + \overline{\beta} \overline{\langle y_2, x \rangle}$$

$$= \overline{\alpha} \langle x, y_1 \rangle + \overline{\beta} \langle x, y_2 \rangle.$$

A vector space V together with an inner product $\langle \cdot, \cdot \rangle$ is termed an inner product space.

Example 16.2. Consider the vector space $V = M_{mn}(\mathbb{F})$. We wish to show that, for some $X, Y \in V$, that $tr(X^*Y)$ defines an inner product in V. Hence we have, $\langle X, Y \rangle \doteq tr(X^*Y) \, \forall X, Y \in V$.

Proof. First suppose that $\langle X, Y \rangle = tr(X^*Y)$ and so;

We check that $tr(X^*Y)$ is Hermitian symmetric:

$$\begin{split} \langle X,Y\rangle &= tr(X^*Y) \\ &= \overline{tr(\overline{X^*Y})} \\ &= \overline{tr(\overline{YX^*})} \\ &= \overline{tr(\overline{(X^*Y)}^T)} \\ &= \overline{tr(Y^*X)} \\ &= \overline{\langle Y,X\rangle}. \end{split}$$

Check that $tr(X^*Y)$ is linear in the first parameters by considering some $\alpha, \beta \in \mathbb{F}$ so that

$$\langle \alpha X_1 + \beta X_2, Y \rangle = tr((\alpha X_1 + \beta X_2)^* Y)$$

$$= tr((\alpha X_1)^* Y + (\beta X_2)^* Y)$$

$$= tr((\alpha X_1)^* Y) + tr((\beta X_2)^* Y)$$

$$= \bar{\alpha} tr(X_1^* Y) + \bar{\beta} tr(X_2^* Y)$$

$$= \bar{\alpha} \langle X_1, Y \rangle + \bar{\beta} \langle X_2, Y \rangle.$$

and check that $tr(X^*Y)$ is positive definite for X = Y,

$$\langle X, X \rangle = tr(X^*X) \ge 0$$

and when

$$X = 0 \implies X^* = 0$$

 $\implies tr(X^*X) = 0.$

TODO..

Example 16.3. Let \mathcal{V} be the vector space of continuous real-valued functions [0,1] endowed with the standard inner product. If f(x) = x and $g(x) = e^x$ then find $\langle f, g \rangle$ and ||f||.

For $\langle f, g \rangle$ we have,

$$\langle f,g\rangle = \int_0^1 f(x)\overline{g(x)}\mathrm{d}x$$

$$= \int_0^1 x e^x \mathrm{d}x$$
 (by parts we have,)
$$= x e^x \|_0^1 - \int_0^1 e^x \mathrm{d}x$$

$$= e - (e - 1) = 1.$$

Now, for ||f|| we have,

$$||f|| = \left(\int_0^1 x^2 dx\right)^{\frac{1}{2}}$$
$$= \sqrt{\frac{1}{x^3}} |_0^1$$
$$= \frac{1}{\sqrt{3}}.$$

Theorem 16.4. Let V be a complex vector space with inner product $\langle \cdot, \cdot \rangle$. Then the inner product induces a norm $\| \cdot \|$ in V by the definition

$$||x|| = \sqrt{\langle x, x \rangle} \text{ for } x \in \mathcal{V}.$$

Proof.

Notice that $\|\cdot\|$ is trivially positive definite as follows:

$$||x|| = \sqrt{\langle x, x \rangle} \ge 0 \Leftrightarrow x = 0, \, \forall x \in \mathcal{V}.$$

By linearity we see that,

$$\begin{split} \|\alpha x\| &= \sqrt{\langle \alpha x, \alpha x \rangle} \\ &= \sqrt{\alpha \bar{\alpha} \langle x, x \rangle} \\ &= \sqrt{|\alpha|^2 \langle x, x \rangle} = |\alpha| \|x\| \, \forall \alpha \in \mathbb{C} \text{ and } \forall x \in \mathcal{V}. \end{split}$$

Finally the triangle inequality follows as shown,

$$||x + y||^{2} = \langle x + y, x + y \rangle$$

$$= ||x||^{2} + ||y||^{2} + \langle x, y \rangle + \langle y, x \rangle$$

$$= ||x||^{2} + ||y||^{2} + \langle x, y \rangle + \overline{\langle x, y \rangle}$$

$$= ||x||^{2} + ||y||^{2} + 2\mathbb{R}[\langle x, y \rangle]$$

$$\leq ||x||^{2} + ||y||^{2} + 2|\langle x, y \rangle|$$

$$\leq ||x||^{2} + ||y||^{2} + 2||x|| ||y||$$

$$= (||x|| + ||y||)^{2} \, \forall x, y \in \mathcal{V}.$$

TODO.. MOVE ME..

Example 16.5. Find a basis for the orthogonal complement of $W = span \left\{ \begin{pmatrix} 1\\1\\1\\0 \end{pmatrix}, \begin{pmatrix} 2\\2\\1\\1 \end{pmatrix} \right\} \in \mathbb{R}^4$. Therefore we do the following:

Therefore we do the following:

$$W^{\perp} = \ker \begin{pmatrix} 1 & 1 & 1 & 0 \\ 2 & 2 & 1 & 1 \end{pmatrix}$$

$$= \ker \begin{pmatrix} 1 & 1 & 1 & 0 \\ 0 & 0 & -1 & 1 \end{pmatrix}$$

$$= span \left\{ \begin{pmatrix} 0 \\ -1 \\ 1 \\ 1 \end{pmatrix}, \begin{pmatrix} -1 \\ 0 \\ 1 \\ 1 \end{pmatrix} \right\}.$$
(after a row reduction)

17. QR DECOMPOSITION

The QR decomposition of a matrix $A \in \mathcal{M}_{n,m}(\mathbb{F})$ into a product of an orthogonal matrix Q and an upper triangular matrix R, written A = QR. However, first we must look at the Gram-Schmidt algorithm in order to obtain an orthonomalised babsis of A.

Definition 17.1 (Projection operator). Consider two $n \times 1$ vectors $\mathbf{u}, \mathbf{v} \in \mathbb{F}^n$. Then the **orthogonal projection** of **v** onto the line spanned by **u** is given by,

$$proj_{\mathbf{u}}(\mathbf{v}) = \frac{\langle \mathbf{v}, \mathbf{u} \rangle}{\langle \mathbf{u}, \mathbf{u} \rangle} \mathbf{u}.$$

Recall that when the field \mathbb{F} is taken to be real for two $n \times 1$ vectors $\mathbf{u}, \mathbf{v} \in \mathbb{R}^n$ then the inner product is then $\langle \mathbf{u}, \mathbf{v} \rangle = \mathbf{u}^T \mathbf{v}$. When the field \mathbb{F} is taken to be complex for two $n \times 1$ vectors $\mathbf{u}, \mathbf{v} \in \mathbb{C}^n$ then the inner product is then $\langle \mathbf{u}, \mathbf{v} \rangle = \mathbf{u}^* \mathbf{v}$. In fact, since the conjugate of a real number is itself the usual complex version can simply be used without issue.

17.1. **Gram-Schmidt.** The Gram-Schmidt algorithm is a method of orthonormalising a finite set of linearly independent vectors that span an inner product space into a orthonormal basis.

Consider a finite sequence of vectors given by the spanning set,

$$S = \{\mathbf{v}_1, \dots, \mathbf{v}_k\}.$$

Then the sequence of Gram-Schmidt orthogonalised vectors, $\{\mathbf{u}_1,\dots,\mathbf{u}_k\}$ is given by,

$$\mathbf{u}_k = \mathbf{v}_k - \sum_{j=1}^{k-1} proj_{\mathbf{u}_k}(\mathbf{v}_k).$$

Hence, we may normalise the sequence by,

$$\mathbf{e}_k = \frac{\mathbf{u}_k}{\|\mathbf{u}_k\|}.$$

Example 17.2.

Consider the spanning set,

$$S = \left\{ \mathbf{v}_1 = \begin{pmatrix} 1 \\ 2 \\ 3 \end{pmatrix}, \mathbf{v}_2 = \begin{pmatrix} 5 \\ 2 \\ 1 \end{pmatrix}, \mathbf{v}_3 = \begin{pmatrix} 3 \\ 2 \\ 4 \end{pmatrix} \right\}.$$

Let $\mathbf{u}_1 = \mathbf{v}_1$ and so,