Math 140C Lecture Notes (Professor: Luca Spolaor)

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## Lecture 1: 4/2/2024

A set  $X \subseteq \mathbb{R}^n$  where  $X \neq \emptyset$  is a vector space if:

- $\overrightarrow{x}, \overrightarrow{y} \in X \Longrightarrow \overrightarrow{x} + \overrightarrow{y} \in X$
- $\vec{x} \in X$  and  $c \in \mathbb{R} \Longrightarrow c\vec{x} \in X$ .

If 
$$\phi = \{\vec{x}_1, \dots, \vec{x}_k\} \subset \mathbb{R}^n$$
, then we define: span  $\phi = \text{span}\{\vec{x}_1, \dots, \vec{x}_k\} = \{c_1\vec{x}_1 + \dots + c_k\vec{x}_k \mid c_1, \dots, c_k \in \mathbb{R}\}.$ 

If  $E \subseteq \mathbb{R}^n$  and  $E = \operatorname{span} \phi$ , then we say  $\phi$  generates E.

Note that  $\mathrm{span}\{\, \overrightarrow{x}_1, \ldots, \, \overrightarrow{x}_2 \}$  forms a vector space (this is trivial to check).

 $\{\vec{x}_1,\ldots,\vec{x}_k\}\subseteq\mathbb{R}^n$  is called linearly independent if:

$$\sum_{i=1}^{k} c_i \vec{x}_i = 0 \Longrightarrow \forall i \in \{1, \dots, k\}, \ c_i = 0.$$

If the above implication does not hold, then we call the set <u>linearly dependent</u>.

If  $X \subseteq \mathbb{R}^n$  is a vector space, then we define the <u>dimension</u> of X as:

$$\dim(X) = \sup\{k \in \mathbb{N} \cup \{0\} \mid \exists \{\vec{x}_1, \dots, \vec{x}_k\} \subset X \text{ which is linearly independent}\}.$$

Also, we define any set containing  $\vec{0}$  to be automatically linearly dependent. This includes the singleton:  $\{\vec{0}\}.$ 

 $Q = \{\overrightarrow{x}_1, \dots, \overrightarrow{x}_k\}$  is a basis for X if:

- ullet Q is linearly independent.
- span Q = X

As an example of a basis, for  $\mathbb{R}^n$  we define the standard basis as the set  $\{e_1, e_2, \dots, e_n\}$  where  $e_i$  is the vector whose ith element is 1 and whose other elements are 0. It is pretty trivial to check that this set is in fact a basis of  $\mathbb{R}^n$ .

<u>Proposition</u>: If  $B = \{\vec{x}_1, \dots, \vec{x}_k\}$  is a basis of a vector space X, then:

**1.** 
$$\forall \vec{v} \in X, c_1, \ldots, c_k \in \mathbb{R} s.t. \vec{v} = \sum_{i=1}^k c_i \vec{x}_i$$

This is true because  $X=\operatorname{span} B.$  So by definition of a span,  $\overrightarrow{v}$  can be expressed as a linear combination of the vectors of B.

2. The  $c_i$  such that  $\overrightarrow{v} = \sum_{i=1}^k c_i \, \overrightarrow{x}_i$  are unique.

Suppose that  $\overrightarrow{v} = \sum c_i \overrightarrow{x}_i = \sum \alpha_i \overrightarrow{x}_i$ . Then  $\overrightarrow{0} = \sum (c_i - \alpha_i) \overrightarrow{x}_i$ . Then since  $\{\overrightarrow{x}_1, \ldots, \overrightarrow{x}_k\}$  are linearly independent, we know for all i that  $c_i - \alpha_i = 0$ . Hence,  $c_i = \alpha_i$  for each i.

Theorem 9.2: Let  $k \in \mathbb{N} \cup \{0\}$ . If  $X = \operatorname{span}\{\overline{x}_1, \dots, \overline{x}_k\}$ , then  $\dim(X) \leq k$ .

Proof:

Suppose for the sake of contradiction that for any  $m \in \mathbb{Z}^+$ , there exists a linearly independent set  $Q = \{ \overrightarrow{y}_1, \dots, \overrightarrow{y}_{k+m} \} \subset X$  which spans X. Then, define  $S_0 = \{ \overrightarrow{x}_1, \dots, \overrightarrow{x}_k \}$  and note that  $S_0$  spans X.

Now by induction, assume for  $i \in \{0,1,\ldots,k-1\}$ , that  $S_i$  contains the first i vectors of Q in addition to k-i vectors of  $S_0$ , and that  $\operatorname{span} S_i = X$ . Then since  $S_i$  spans X, we know that  $y_{i+1} \in X$  is in the span of  $S_i$ . So, letting  $\overrightarrow{x}_{n_1},\ldots,\overrightarrow{x}_{n_{k-i}}$  be the elements from  $S_0$  in  $S_i$ , we know that there exists scalars  $a_1,\ldots,a_{i+1},b_1,\ldots,b_{k-i}\in\mathbb{R}$  where  $a_{i+1}=1$  such that:

$$\sum_{j=1}^{i+1} a_j \, \overrightarrow{y}_j + \sum_{j=1}^{k-i} b_j \, \overrightarrow{x}_{n_j} = \, \overrightarrow{0}$$

If all  $b_j=0$ , then we have a contradiction. This is because  $\{\vec{y}_1,\ldots,\vec{y}_{k+1}\}$  is assumed to be linearly independent. So, having all  $b_j=0$  implies that:

$$\sum_{j=1}^{i+1} a_j \, \vec{y}_j = \sum_{j=1}^{i+1} a_j \, \vec{y}_j + \sum_{j=i+2}^{k+1} 0 \cdot \, \vec{y}_j = \, \vec{0}$$

In turn this means that all  $a_j=0$ , which contradicts that  $a_{i+1}=1$ .

So, not all  $b_j=0$ . This means that for some j we must have that  $\overrightarrow{x}_{n_j}$  is in the span of  $(S_i\setminus\{\overrightarrow{x}_{n_j}\})\cup\{\overrightarrow{y}_{i+1}\}$ . Call this set  $S_{i+1}$ . Clearly,  $S_{i+1}$  contains the first i+1 vectors of Q. Also:

$$\operatorname{span} S_{i+1} = \operatorname{span} (S_i \cup \{ \overrightarrow{y}_{i+1} \}) = \operatorname{span} S_i = X.$$

So  $S_{i+1}$  satisfies the same conditions  $S_i$  did.

Now we get to the contradiction. Using the above reasoning, we will eventually construct  $S_k = \{ \overrightarrow{y}_1, \dots, \overrightarrow{y}_k \}$  which still spans X. However, since  $\overrightarrow{y}_{k+1} \in X$ , that means that  $\overrightarrow{y}_{k+1}$  equals some linear combination of the other  $\overrightarrow{y}$  in Q. This contradicts that Q is linearly independent.  $\blacksquare$ 

Corollary: If  $B = \{ \overrightarrow{x}_1, \dots, \overrightarrow{x}_k \}$  is a basis for X, then  $\dim(X) = k$ .

Proof:

Since B is linearly independent, by definition  $\dim(X) \geq k$ . Meanwhile, since B spans X, we know by the above theorem that  $\dim(X) \leq k$ . So  $\dim(X) = k$ .

Theorem 9.3: Suppose X is a vector space and dim(X) = n. Then:

(A) For  $E = \{\vec{x}_1, \dots, \vec{x}_n\} \subset X$ , we have that  $X = \operatorname{span} E$  if and only if E is linearly independent.

Proof:

First, assume E is linearly independent. Then, note that for any  $\overrightarrow{y} \in X$ , we must have that  $E \cup \{\overrightarrow{y}\}$  is linearly dependent because  $|E \cup \{\overrightarrow{y}\}| > \dim(X)$ . So, there exists  $c_1, \ldots, c_n, c_{n+1} \in \mathbb{R}$  such that at least one  $c_i$  is nonzero and:

$$\sum_{i=1}^{n} c_i \, \overrightarrow{x}_i + c_{n+1} \, \overrightarrow{y} = \, \overrightarrow{0}$$

Now if  $c_{n+1}=0$ , we have a contradiction because E is linearly independent. So, we conclude that  $c_{n+1}\neq 0$ . Thus, by rearranging terms we can express y as a linear combination of the vectors of E. Therefore,  $\operatorname{span} E=X$  since y can be any vector in X.

Secondly, assume E is not linearly independent. Then for some  $\overrightarrow{x}_i \in E$ , we have that  $\operatorname{span} E = \operatorname{span}(E \setminus \{\overrightarrow{x}_i\})$ . However,  $|E \setminus \{\overrightarrow{x}_i\}| = n-1$ . So if  $X = \operatorname{span} E$ , then  $\dim(X) \leq |E \setminus \{\overrightarrow{x}_i\}| = n-1$ , which contradicts our assumption that  $\dim(X) = n$ . Hence,  $X \neq \operatorname{span} E$ .

(B) X has a basis and every basis of X consists of n vectors.

Proof:

By the definition of  $\dim(X)$ , we know that there exists a linearly independent set of n vectors. By the previous part of this theorem, we also know that that set spans X. So, it is a basis of X. Meanwhile, by the corollary to theorem 9.2, we know that the number of vectors in a basis of X equals the dimension of X. Hence, all bases of X must have n vectors.

(C) If  $1 \leq m \leq n$  and  $\{\overrightarrow{y}_1, \ldots, \overrightarrow{y}_m\} \subset X$  is linearly independent, then X has a basis that contains  $\overrightarrow{y}_1, \ldots, \overrightarrow{y}_m$ .

Proof:

Let  $S_0 = \{\vec{x}_1, \dots, \vec{x}_n\}$  be a basis of X and  $Q = \{\vec{y}_1, \dots, \vec{y}_m\}$ . Then by the same induction which we used to prove theorem 9.2, we can construct a basis:  $S_m$ , of X which contains  $\vec{y}_1, \dots, \vec{y}_m$ .

Let X and Y be vector spaces. A map  $\mathbf{A}: X \longrightarrow Y$  is <u>linear</u> if  $\mathbf{A}(c_1 \overrightarrow{x}_1 + c_2 \overrightarrow{x}_2) = c_1 \mathbf{A}(\overrightarrow{x}_1) + c_2 \mathbf{A}(\overrightarrow{x}_2)$  for all  $\overrightarrow{x}_1, \overrightarrow{x}_2 \in X$  and  $c_1, c_2 \in \mathbb{R}$ .

**Observations:** 

1. A linear map sends  $\vec{0}$  to  $\vec{0}$ . This is because:

$$\mathbf{A}(\vec{0}) = \mathbf{A}(\vec{v} - \vec{v}) = \mathbf{A}(\vec{v}) - \mathbf{A}(\vec{v}) = \vec{0}.$$

2. If  $\mathbf{A}: X \longrightarrow Y$  is a linear map and  $B = \{\overrightarrow{x}_1, \dots, \overrightarrow{x}_k\}$  is a basis of X, then  $\mathbf{A}\left(\sum\limits_{i=1}^k (c_i \overrightarrow{x}_i)\right) = \sum\limits_{i=1}^k c_i \mathbf{A}(\overrightarrow{x}_i)$  for all  $c_1, \dots, c_k \in \mathbb{R}$ .

Given two vector spaces X and Y, we define L(X,Y) to be the set of all linear transformations from X into Y. Also, we shall abbreviate L(X,X) as L(X).

$$\mathcal{N}(\mathbf{A}) = \text{"null space / kernel of } \mathbf{A} \text{"} = \{ \overrightarrow{x} \in X \mid \mathbf{A}(\overrightarrow{x}) = \overrightarrow{0} \}.$$

$$\mathscr{R}(\mathbf{A}) = \text{"range of } \mathbf{A} \text{"} = \{ \overrightarrow{y} \in Y \mid \exists \overrightarrow{x} \in X \ s.t. \ \mathbf{A} \overrightarrow{x} = \overrightarrow{y} \}.$$

<u>Proposition</u>: For any linear map  $\mathbf{A}: X \longrightarrow Y$ ,  $\mathcal{N}(\mathbf{A})$  and  $\mathcal{R}(\mathbf{A})$  are vector spaces.

Proof:

- Assume  $\vec{x}_1, \vec{x}_2 \in \mathcal{N}(\mathbf{A}) \subset X$  and  $c \in \mathbb{R}$ . Then:
  - $\mathbf{A}(\overrightarrow{x}_1 + \overrightarrow{x}_2) = \mathbf{A}(\overrightarrow{x}_1) + \mathbf{A}(\overrightarrow{x}_2) = \overrightarrow{0} + \overrightarrow{0} = \overrightarrow{0}$ , which means that  $\overrightarrow{x}_1 + \overrightarrow{x}_2 \in \mathcal{N}(\mathbf{A})$ .
  - $\circ \mathbf{A}(c\overrightarrow{x}_1) = c\mathbf{A}(\overrightarrow{x}_1) = c\overrightarrow{0} = \overrightarrow{0}. \text{ So } c\overrightarrow{x}_1 \in \mathscr{N}(\mathbf{A}).$

This shows that  $\mathcal{N}(\mathbf{A})$  is a vector space.

- Assume  $\vec{y}_1, \vec{y}_2 \in \mathcal{R}(\mathbf{A}) \subset Y$  and  $c \in \mathbb{R}$ . Then:
  - $\begin{array}{l} \circ \ \ \text{We know there exists} \ \overrightarrow{x}_1, \ \overrightarrow{x}_2 \in X \ \text{such that} \ \mathbf{A}(\overrightarrow{x}_1) = \ \overrightarrow{y}_1 \ \text{and} \\ \mathbf{A}(\overrightarrow{x}_2) = \ \overrightarrow{y}_2. \ \text{In turn,} \ \mathbf{A}(\overrightarrow{x}_1 + \overrightarrow{x}_2) = \mathbf{A}(\overrightarrow{x}_1) + \mathbf{A}(\overrightarrow{x}_2) = \ \overrightarrow{y}_1 + \ \overrightarrow{y}_2. \\ \text{So} \ \overrightarrow{y}_1 + \ \overrightarrow{y}_2 \in \mathscr{R}(\mathbf{A}). \end{array}$
  - $\circ$  Now continue letting  $\overrightarrow{x}_1 \in X$  be a vector such that  $\mathbf{A}(\overrightarrow{x}_1) = \overrightarrow{y}_1$ . Then  $\mathbf{A}(c\overrightarrow{x}_1) = c\mathbf{A}(\overrightarrow{x}_1) = c\overrightarrow{y}_1$ . So  $c\overrightarrow{y}_1 \in \mathscr{R}(\mathbf{A})$ .

This shows that  $\mathcal{R}(\mathbf{A})$  is a vector space.

$$\operatorname{rk}(\mathbf{A}) = \text{"rank of } \mathbf{A} \text{"} = \dim(\mathscr{R}(\mathbf{A})).$$

$$\operatorname{null}(\mathbf{A}) = "\underline{\operatorname{nullity}} \text{ of } \mathbf{A}" = \dim(\mathscr{N}(\mathbf{A})).$$

Rank-Nullity Theorem: Given any  $\mathbf{A} \in L(X,Y)$ , we have that  $\dim(X) = \mathrm{rk}(\mathbf{A}) + \mathrm{null}(\mathbf{A})$ .

Proof:

Let 
$$\dim(X) = n$$
.

 $\mathscr{N}(\mathbf{A})\subseteq X$  is a vector space. So pick a basis  $\{\overrightarrow{v}_1,\ldots,\overrightarrow{v}_k\}$  for  $\mathscr{N}(\mathbf{A})$  where  $k=\mathrm{null}(\mathbf{A})\leq \dim(X)$ . Then by theorem 9.3, choose  $\overrightarrow{w}_1,\ldots,\overrightarrow{w}_{m-k}$  such that  $\{\overrightarrow{v}_1,\ldots,\overrightarrow{v}_k,\overrightarrow{w}_1,\ldots,\overrightarrow{w}_{n-k}\}$  is a basis of X. Note that  $\dim(X)=n$ .

Claim:  $B = {\mathbf{A}(\vec{w}_1), \dots, \mathbf{A}(\vec{w}_{n-k})}$  is a basis of  $\mathcal{R}(\mathbf{A})$ .

•  $\mathbf{A}(\overrightarrow{v_i}) = \overrightarrow{0}$  for all  $i \in \{1, \dots, k\}$ . So:

$$\mathcal{R}(\mathbf{A}) = \operatorname{span}\{\mathbf{A}(\overrightarrow{v}_1), \dots, \mathbf{A}(\overrightarrow{v}_k), \mathbf{A}(\overrightarrow{w}_1), \dots, \mathbf{A}(\overrightarrow{w}_{n-k})\}$$
$$= \operatorname{span}\{\mathbf{A}(\overrightarrow{w}_1), \dots, \mathbf{A}(\overrightarrow{w}_{n-k})\} = \operatorname{span} B$$

ullet B is linearly independent.

To see this, note that: 
$$\sum_{i=1}^{n-k} (c_i \mathbf{A}(\overrightarrow{w}_i)) = \overrightarrow{0} \Longrightarrow \mathbf{A} \left( \sum_{i=1}^{n-k} c_i \overrightarrow{w}_i \right) = \overrightarrow{0}$$

Since we picked each  $\overrightarrow{w}_1,\ldots,\overrightarrow{w}_{n-k}\in B$  so that they were not in  $\mathcal{N}(A)$ , we know that any vector in the span of B is not mapped to 0 by  $\mathbf{A}$  unless it is the zero vector. So

$$\sum_{i=1}^{n-k} c_i \vec{w}_i = \vec{0}$$

And since all the  $\vec{w}_i$  are linearly independent, all constants  $c_i$  equal 0.

So 
$$\operatorname{rk}(\mathbf{A}) = n - k = \dim(X) - \operatorname{null}(\mathbf{A}).$$

## Lecture 2: 4/4/2024

<u>Proposition</u>: Given  $A \in L(X, Y)$ , then:

• **A** is injective if and only if  $null(\mathbf{A}) = 0$ .

Proof:

( $\Longrightarrow$ ) If  ${\bf A}$  is injective, then since  ${\bf A}(\vec{0})=\vec{0}$ , we have that any vector  $\vec{v}\neq\vec{0}$  is not in  $\mathscr{N}({\bf A})$ . So  $\mathscr{N}({\bf A})=\{\vec{0}\}$ , meaning  $\operatorname{null}({\bf A})=0$ .

( $\iff$ ) If  $\operatorname{null}(\mathbf{A})=0$ , then  $\mathbf{A}(\overrightarrow{v})=\overrightarrow{0} \implies \overrightarrow{v}=\overrightarrow{0}$ . So now assume  $\mathbf{A}(\overrightarrow{v})=\mathbf{A}(\overrightarrow{u})$ . Then  $\mathbf{A}(\overrightarrow{v}-\overrightarrow{u})=\overrightarrow{0}$ , meaning  $\overrightarrow{v}=\overrightarrow{u}$ . Hence  $\mathbf{A}$  is injective.

• **A** is surjective if and only if  $rk(\mathbf{A}) = dim(Y)$ .

Proof:

( $\Longrightarrow$ ) If **A** is surjective then  $\mathcal{R}(\mathbf{A})=Y$ . So we automatically have that  $\mathrm{rk}(\mathbf{A})=\dim(Y)$ 

( $\Leftarrow$ ) If  $\operatorname{rk}(\mathbf{A}) = \dim(Y)$ , then there exists a linearly independent set of vectors  $B \subset \mathscr{R}(\mathbf{A})$  containing  $\dim(Y)$  many vectors and spanning  $\mathscr{R}(\mathbf{A})$ . Then by theorem 9.3, since  $B \subset \mathscr{R}(\mathbf{A}) \subseteq Y$ , we know  $\operatorname{span} B = Y$ . So,  $\mathscr{R}(\mathbf{A}) = Y$ , meaning  $\mathbf{A}$  is surjective.

<u>Corollary</u>: Let  $A \in L(X)$ . Then A is bijective if and only if null(A) = 0.

Proof: (let  $A: X \longrightarrow X$  be a linear map)

 $(\Longrightarrow)$  If  ${\bf A}$  is bijective, then automatically  ${\bf A}$  is injective. So  ${\rm null}({\bf A})=0$  by the previous proposition.

( $\Leftarrow$ ) If  $\operatorname{null}(\mathbf{A}) = 0$ , then by the rank-nullity theorem, we know that  $\operatorname{rk}(\mathbf{A}) = \dim(X)$ . Thus  $\mathbf{A}$  is both injective and surjective, meaning  $\mathbf{A}$  is bijective.

For  $\mathbf{A} \in L(X)$ , when  $\operatorname{null}(\mathbf{A}) = 0$ , we call  $\mathbf{A}$  invertible and define  $\mathbf{A}^{-1} : X \longrightarrow X$  by  $\mathbf{A}^{-1}(\mathbf{A}(\overrightarrow{x})) = \overrightarrow{x}$  for all  $\overrightarrow{x} \in X$ .

Because **A** must be a bijective set function, we know that  $\mathbf{A}^{-1}$  must also be a right-inverse of **A**, meaning  $\mathbf{A}(\mathbf{A}^{-1}(\vec{x})) = \vec{x}$ .

Additionally, consider any  $\vec{x}_1, \vec{x}_2 \in X$  and let  $\vec{x}_1' = \mathbf{A}^{-1}(\vec{x}_1)$  and  $\vec{x}_2' = \mathbf{A}^{-1}(\vec{x}_2)$ . Then since  $\mathbf{A}$  is a linear mapping, we know that for any  $c_1, c_2 \in \mathbb{R}$ :

$$\mathbf{A}(c_1 \, \vec{x}_1' + c_2 \, \vec{x}_2') = c_1 \mathbf{A}(\mathbf{A}^{-1}(\, \vec{x}_1)) + c_2 \mathbf{A}(\mathbf{A}^{-1}(\, \vec{x}_2)) = c_1 \, \vec{x}_1 + c_2 \, \vec{x}_2$$

So:  $\mathbf{A}^{-1}(c_1\vec{x}_1 + c_2\vec{x}_2) = c_1\vec{x}_1' + c_2\vec{x}_2' = c_1\mathbf{A}^{-1}(\vec{x}_1) + c_2\mathbf{A}^{-1}(\vec{x}_2)$ . Hence, we've shown that  $\mathbf{A}^{-1}$  is a linear mapping, meaning that  $\mathbf{A}^{-1} \in L(X)$ .

Let  $\mathbf{A} \in L(X,Y)$  and  $\mathbf{B} \in L(Y,Z)$ . Then we define  $\mathbf{B}\mathbf{A}: X \longrightarrow Z$  by the rule that  $\overrightarrow{x} \mapsto \mathbf{B}(\mathbf{A}(\overrightarrow{x}))$ .

We can trivially show that **BA** is a linear mapping. Consider any  $\vec{x}_1, \vec{x}_2 \in X$  and  $c_1, c_2 \in \mathbb{R}$ . Then:

$$\mathbf{B}\mathbf{A}(c_1 \, \overrightarrow{x}_1 + c_2 \, \overrightarrow{x}_2) = \mathbf{B}(c_1 \mathbf{A}(\, \overrightarrow{x}_1) + c_2 \mathbf{A}(\, \overrightarrow{x}_2))$$

$$= c_1 \mathbf{B}(\mathbf{A}(\, \overrightarrow{x}_1)) + c_2 \mathbf{B}(\mathbf{A}(\, \overrightarrow{x}_2))$$

$$= c_1 \mathbf{B}\mathbf{A}(\, \overrightarrow{x}_1) + c_2 \mathbf{B}\mathbf{A}(\, \overrightarrow{x}_2)$$

This means that  $\mathbf{BA} \in L(X, Z)$ .

Let  $\mathbf{A}, \mathbf{B} \in L(X, Y)$  and  $c_1, c_2 \in \mathbb{R}$ . Then we define  $(c_1\mathbf{A} + c_2\mathbf{B}) : X \longrightarrow Y$  by the rule:  $\overrightarrow{x} \mapsto c_1\mathbf{A}(\overrightarrow{x}) + c_2\mathbf{B}(\overrightarrow{x})$ .

It is even more trivial to show that  $(c_1\mathbf{A} + c_2\mathbf{B})$  is a linear map.

Let  $\mathbf{A} \in L(\mathbb{R}^n, \mathbb{R}^m)$ . We define the <u>norm</u> of  $\mathbf{A}$  as:  $\|\mathbf{A}\| = \sup \{\|\mathbf{A}(\overrightarrow{x})\| \mid \overrightarrow{x} \in \mathbb{R}^n \text{ and } \|\overrightarrow{x}\| \leq 1\}.$ 

Throughout this section, we shall prove that  $\|\cdot\|:L(\mathbb{R}^n,\mathbb{R}^m)\longrightarrow\mathbb{R}$  is well-defined and fulfills the properties of a general norm function.

<u>Proposition</u>: If  $A \in L(\mathbb{R}^n, \mathbb{R}^m)$ , then ||A|| exists and is finite.

Proof:

Let  $\{e_1, \ldots, e_n\}$  be the standard basis in  $\mathbb{R}^n$ . Then for any  $\vec{x} \in \mathbb{R}^n$ , there are unique  $c_1, \ldots, c_n \in \mathbb{R}$  such that  $\vec{x} = c_1 e_1 + \ldots + c_n e_n$ .

Since we are working with the standard basis, we know:  $\|\vec{x}\| = \sqrt{\sum_{i=1}^{n} c_i^2}$ .

Thus, for  $\|\vec{x}\| \le 1$ , we must have that  $|c_i| \le 1$  for each  $c_i$ . This means:

$$\|\mathbf{A}(\vec{x})\| = \left\|\sum_{i=1}^{n} c_i \mathbf{A}(e_i)\right\| \le \sum_{i=1}^{n} \|c_i \mathbf{A}(e_i)\| = \sum_{i=1}^{n} |c_i| \|\mathbf{A}(e_i)\| \le \sum_{i=1}^{n} \|\mathbf{A}(e_i)\|$$

Importantly, we must have that  $\sum_{i=1}^{n} \|\mathbf{A}(e_i)\|$  is finite. Additionally, it is an upper bound to the set:  $\{\|\mathbf{A}(\overrightarrow{x})\| \mid \overrightarrow{x} \in \mathbb{R}^n \text{ and } \|\overrightarrow{x}\| \leq 1\} \subseteq \mathbb{R}$ .

So, we showed that the above set is bounded above. Also, the above set is nonempty because it must contain  $\|\vec{0}\| = 0$ . Thus by the least upper bound property of  $\mathbb{R}$ , we know that the supremum of this set exists in  $\mathbb{R}$ .

Hence,  $\|\mathbf{A}\|$  exists and is finite.

Side note, the above proof also shows that  $\|\mathbf{A}\| \geq 0$ .

 $\underline{\text{Lemma}} \text{: For } \mathbf{A} \in L(\mathbb{R}^n, \mathbb{R}^m) \text{ and } \overrightarrow{x} \in \mathbb{R}^n \text{, we have that } \|\mathbf{A}(\overrightarrow{x})\| \leq \|\mathbf{A}\| \|\overrightarrow{x}\|.$ 

Proof:

Case 1:  $\vec{x} \neq \vec{0}$ .

Then since  $\|\vec{x}\| \neq 0$ , we can say that:

$$\|\mathbf{A}(\vec{x})\| = \|\mathbf{A}(\vec{x})\| \frac{\vec{x}}{\|\vec{x}\|} \|\mathbf{A}(\vec{x})\| \frac{\vec{x}}{\|\vec{x}\|} \|\mathbf{A}(\vec{x})\| = \|\mathbf{A}(\vec{x})\| \|\mathbf{A}($$

Now 
$$\frac{\overrightarrow{x}}{\|\overrightarrow{x}\|} \in \mathbb{R}^n$$
 and  $\left\|\frac{\overrightarrow{x}}{\|\overrightarrow{x}\|}\right\| = 1$ . So,  $\left\|\mathbf{A}\left(\frac{\overrightarrow{x}}{\|\overrightarrow{x}\|}\right)\right\| \|\overrightarrow{x}\| \le \|\mathbf{A}\| \|\overrightarrow{x}\|$ 

Case 2: 
$$\vec{x} = \vec{0}$$
.

Then trivially 
$$\|\mathbf{A}(\overrightarrow{x})\| = \|\mathbf{A}(\overrightarrow{0})\| = 0 = \|\mathbf{A}\| \|\overrightarrow{0}\| = \|\mathbf{A}\| \|\overrightarrow{x}\|$$

<u>Proposition</u>: If  $\mathbf{A} \in L(\mathbb{R}^n, \mathbb{R}^m)$ , then  $0 \le \|\mathbf{A}\|$ . Also  $\|\mathbf{A}\| = 0$  if and only if  $\mathbf{A}$  is the unique function mapping all of  $\mathbb{R}^n$  to  $\overrightarrow{0}$ .

Proof:

We already showed previously that  $\|\mathbf{A}\| \geq 0$ . So, it now suffices to show that  $\|\mathbf{A}\| = 0 \iff \mathcal{N}(\mathbf{A}) = \mathbb{R}^n$ .

( $\Longrightarrow$ ) Assume that  $\mathcal{N}(\mathbf{A}) \neq \mathbb{R}^n$ . Then there exists  $\overrightarrow{x} \in \mathbb{R}^n$  such that  $\mathbf{A}(\overrightarrow{x}) \neq \overrightarrow{0}$ . Since  $\overrightarrow{x}$  can't be  $\overrightarrow{0}$ , consider the vector  $\hat{x} = \frac{\overrightarrow{x}}{\|\overrightarrow{x}\|}$ . By the linearity of  $\mathbf{A}$ , we know  $\mathbf{A}(\hat{x}) = \frac{1}{\|\overrightarrow{x}\|}\mathbf{A}(\overrightarrow{x}) \neq \overrightarrow{0}$ . So,  $\|\mathbf{A}(\hat{x})\| > 0$ . But  $\|\mathbf{A}(\hat{x})\|$  is in the set that  $\|\mathbf{A}\|$  is a supremum of, which means that  $\|\mathbf{A}\| \geq \|\mathbf{A}(\hat{x})\| > 0$ . Or in other words,  $\|\mathbf{A}\| \neq 0$ .

(
$$\Leftarrow$$
) Assume that  $\mathcal{N}(\mathbf{A})=\mathbb{R}^n$ . Then,  $\sup\{\|\mathbf{A}(\overrightarrow{x})\|\mid \overrightarrow{x}\in\mathbb{R}^n \text{ and } \|\overrightarrow{x}\|\leq 1\}=\sup\{0\}=0$ 

<u>Corollary</u>: Given  $\mathbf{A} \in L(\mathbb{R}^n, \mathbb{R}^m)$ , we have that  $\mathbf{A}$  is uniformly continuous.

Proof:

Case 1:  $\|\mathbf{A}\| \neq 0$ , meaning we can divide by  $\|\mathbf{A}\|$ . By the previous proposition,  $\|\mathbf{A}(\overrightarrow{x}) - \mathbf{A}(\overrightarrow{y})\| \leq \|\mathbf{A}\| \|\overrightarrow{x} - \overrightarrow{y}\|$  for all  $\overrightarrow{x}, \overrightarrow{y} \in \mathbb{R}^n$ . Hence, for any  $\varepsilon > 0$ , if we make  $\|\overrightarrow{x} - \overrightarrow{y}\| < \frac{\varepsilon}{\|\mathbf{A}\|}$ , then  $\|\mathbf{A}(\overrightarrow{x}) - \mathbf{A}(\overrightarrow{y})\| < \varepsilon$ .

Case 2:  $\|\mathbf{A}\| = 0$ .

Then  ${\bf A}$  is a constant function, making it automatically uniformly continuous.

<u>Subcorollary</u>: Given  $\mathbf{A} \in L(\mathbb{R}^n, \mathbb{R}^m)$ . there exists  $\vec{x} \in \mathbb{R}^n$  with  $\|\vec{x}\| \leq 1$  such that  $\|\mathbf{A}(\vec{x})\| = \|\mathbf{A}\|$ .

Proof:

Let  $S = \{ \overrightarrow{x} \in \mathbb{R}^n \mid ||\overrightarrow{x}|| \le 1 \}$  and consider the restriction  $\mathbf{A}|_S$ .

Since S is a closed and bounded subset of  $\mathbb{R}^n$ , we know that S is compact by the Heine-Borel theorem (see proposition 28 in Math 140A notes). This combined with the fact that  $\mathbf{A}|_S$  is still continuous means that by the extreme value theorem, there is  $\overrightarrow{x} \in S$  with:

$$\mathbf{A}(\overrightarrow{x}) = \mathbf{A}|_{S}(\overrightarrow{x}) = \sup \{ \|\mathbf{A}(\overrightarrow{x})\| \mid \overrightarrow{x} \in \mathbb{R}^{n} \text{ and } \|\overrightarrow{x}\| \leq 1 \}.$$

<u>Proposition</u>: If  $A, B \in L(\mathbb{R}^n, \mathbb{R}^m)$ , then  $||A + B|| \le ||A|| + ||B||$ .

Proof:

Let 
$$\overrightarrow{x} \in \mathbb{R}^n$$
 be a vector such that  $\|\overrightarrow{x}\| \le 1$  and  $\|\mathbf{A}(\overrightarrow{x})\| = \|\mathbf{A}\|$ . Then:  $\|\mathbf{A} + \mathbf{B}\| = \|(\mathbf{A} + \mathbf{B})(\overrightarrow{x})\| = \|\mathbf{A}(\overrightarrow{x}) + \mathbf{B}(\overrightarrow{x})\|$   $\leq \|\mathbf{A}(\overrightarrow{x})\| + \|\mathbf{B}(\overrightarrow{x})\| \leq \|\mathbf{A}\| + \|\mathbf{B}\|$ 

<u>Proposition</u>: If  $\mathbf{A} \in L(\mathbb{R}^n, \mathbb{R}^m)$  and  $c \in \mathbb{R}$ , then  $||c\mathbf{A}|| = |c|||\mathbf{A}||$ .

Proof:

Pick 
$$\overrightarrow{x} \in \mathbb{R}^n$$
 satisfying  $\|\overrightarrow{x}\| \le 1$  and  $\|\mathbf{A}(\overrightarrow{x})\| = \|\mathbf{A}\|$ . Then:  $|c|\|\mathbf{A}\| = |c|\|\mathbf{A}(\overrightarrow{x})\| = \|c\mathbf{A}(\overrightarrow{x})\| = \|(c\mathbf{A})(\overrightarrow{x})\| \le \|c\mathbf{A}\|$ .

Next, pick 
$$\overrightarrow{y} \in \mathbb{R}^n$$
 satisfying  $\|\overrightarrow{y}\| \le 1$  and  $\|(c\mathbf{A})(\overrightarrow{x})\| = \|c\mathbf{A}\|$ . Then:  $\|c\mathbf{A}\| = \|(c\mathbf{A})(\overrightarrow{y})\| = \|c\mathbf{A}(\overrightarrow{y})\| = |c|\|\mathbf{A}\overrightarrow{y}\| \le |c|\|\mathbf{A}\|$ .

Specifically because of the four propositions above, we have shown that  $\|\cdot\|:L(\mathbb{R}^n,\mathbb{R}^m)\longrightarrow\mathbb{R}$  is well-defined and a valid norm. Consequently, by defining  $d(\mathbf{A},\mathbf{B})=\|\mathbf{A}-\mathbf{B}\|$  for all  $\mathbf{A},\mathbf{B}\in L(\mathbb{R}^n,\mathbb{R}^m)$ , we naturally get that  $L(\mathbb{R}^n,\mathbb{R}^m)$  is a metric space.

Given any  $A, B, C \in L(\mathbb{R}^n, \mathbb{R}^m)$ , we have:

- $d(\mathbf{A}, \mathbf{B}) = \|\mathbf{A} \mathbf{B}\| \ge 0$  with  $d(\mathbf{A}, \mathbf{B}) = 0$ Also  $d(\mathbf{A}, \mathbf{B}) = 0$  if and only if  $\mathbf{A} = \mathbf{B}$ .
- $d(\mathbf{A}, \mathbf{B}) = ||\mathbf{A} \mathbf{B}|| = |-1|||\mathbf{B} \mathbf{A}|| = d(\mathbf{B}, \mathbf{A})$
- $d(\mathbf{A}, \mathbf{C}) = \|\mathbf{A} \mathbf{C}\| \le \|\mathbf{A} \mathbf{B}\| + \|\mathbf{B} \mathbf{C}\| = d(\mathbf{A}, \mathbf{B}) + d(\mathbf{B}, \mathbf{C})$

Before moving on, here is another corollary of the above statements.

Corollary: If 
$$\mathbf{A} \in L(\mathbb{R}^n, \mathbb{R}^m)$$
 and  $\mathbf{B} \in L(\mathbb{R}^m, \mathbb{R}^k)$ , then  $\|\mathbf{B}\mathbf{A}\| \leq \|\mathbf{B}\| \|\mathbf{A}\|$ .

Proof:

Pick 
$$\overrightarrow{x} \in \mathbb{R}^n$$
 satisfying  $\|\overrightarrow{x}\| \le 1$  and  $\|(\mathbf{B}\mathbf{A})(\overrightarrow{x})\| = \|\mathbf{B}\mathbf{A}\|$ . Then:  $\|\mathbf{B}\mathbf{A}\| = \|(\mathbf{B}\mathbf{A})(\overrightarrow{x})\| = \|\mathbf{B}(\mathbf{A}(\overrightarrow{x}))\| \le \|\mathbf{B}\|\|\mathbf{A}(\overrightarrow{x})\| \le \|\mathbf{B}\|\|\mathbf{A}\|$ .

<u>Theorem 9.8</u>: Let  $\Omega \subset L(\mathbb{R}^n)$  be the set of all invertible linear mappings on  $\mathbb{R}^n$ .

(A) If 
$$\mathbf{A}\in\Omega$$
,  $\mathbf{B}\in L(\mathbb{R}^n)$ , and  $\|\mathbf{B}-\mathbf{A}\|<\frac{1}{\|\mathbf{A}^{-1}\|}$ , then  $\mathbf{B}\in\Omega$ .

Proof:

Pick  $\overrightarrow{x} \in \mathbb{R}^n$  such that  $\|\overrightarrow{x}\| \leq 1$ . Then:

$$\begin{split} \|\mathbf{A}(\overrightarrow{x})\| &= \|(\mathbf{A} - \mathbf{B} + \mathbf{B})(\overrightarrow{x})\| \\ &\leq \|(\mathbf{A} - \mathbf{B})(\overrightarrow{x})\| + \|\mathbf{B}(\overrightarrow{x})\| \\ &\leq \|\mathbf{A} - \mathbf{B}\| \|\overrightarrow{x}\| + \|\mathbf{B}(\overrightarrow{x})\| = \|\mathbf{B} - \mathbf{A}\| \|\overrightarrow{x}\| + \|\mathbf{B}(\overrightarrow{x})\| \end{split}$$

Meanwhile, note that  $\|\mathbf{A}^{-1}\| \neq 0$ . We know this because  $\mathbf{A}^{-1}$  must be invertible (because  $\mathcal{N}(\mathbf{A}^{-1}) = \{ \overrightarrow{0} \}$ ) and the one linear transformation in  $L(\mathbb{R}^n)$  with norm 0 is not invertible. So:

$$\tfrac{\parallel\overrightarrow{x}\parallel}{\parallel\mathbf{A}^{-1}\parallel}=\tfrac{\parallel\mathbf{A}^{-1}\mathbf{A}(\overrightarrow{x})\parallel}{\parallel\mathbf{A}^{-1}\parallel}\leq \tfrac{\parallel\mathbf{A}^{-1}\parallel\parallel\mathbf{A}(\overrightarrow{x})\parallel}{\parallel\mathbf{A}^{-1}\parallel}=\|\mathbf{A}(\overrightarrow{x})\|$$

Hence,  $\frac{\|\overrightarrow{x}\|}{\|\mathbf{A}^{-1}\|} \leq \|\mathbf{B} - \mathbf{A}\| \|\overrightarrow{x}\| + \|\mathbf{B}(\overrightarrow{x})\|$ . By rearranging terms, we get this expression:  $\left(\frac{1}{\|\mathbf{A}^{-1}\|} - \|\mathbf{B} - \mathbf{A}\|\right) \|\overrightarrow{x}\| \leq \|\mathbf{B}(\overrightarrow{x})\|$ .

Now, note that if  $\|\mathbf{B}(\overrightarrow{x})\| = 0$  but  $\overrightarrow{x} \neq \overrightarrow{0}$ , then we must have that:  $\frac{1}{\|\mathbf{A}^{-1}\|} - \|\mathbf{B} - \mathbf{A}\| \leq 0$ . Or in other words,  $\|\mathbf{B} - \mathbf{A}\| \geq \frac{1}{\|\mathbf{A}^{-1}\|}$ . So, if  $\|\mathbf{B} - \mathbf{A}\| < \frac{1}{\|\mathbf{A}^{-1}\|}$ , then  $\|\mathbf{B}(\overrightarrow{x})\| = 0$  only when  $\overrightarrow{x} = \overrightarrow{0}$ . Hence,  $\mathrm{null}(\mathbf{B}) = 0$  and  $\mathbf{B}$  is invertible.

(B)  $\Omega$  is an open subset of  $L(\mathbb{R}^n)$ , and the mapping over  $\Omega$  with the rule:  $\mathbf{A}\mapsto \mathbf{A}^{-1}$ , is continuous.

Proof:

Firstly, by part A we know that for any  $\mathbf{A} \in \Omega$ , if  $r = \frac{1}{\|\mathbf{A}^{-1}\|}$ , then  $B_r(\mathbf{A}) \subseteq \Omega$ . So,  $\Omega$  is an open set in the metric space  $L(\mathbb{R}^n)$ .

Now let  $A, B \in \Omega$  and recall from part A that:

$$\left(\frac{1}{\|\mathbf{A}^{-1}\|} - \|\mathbf{B} - \mathbf{A}\|\right) \|\overrightarrow{x}\| \le \|\mathbf{B}(\overrightarrow{x})\|.$$

Since we know  $\mathbf{B}^{-1}$  exists, set  $\overrightarrow{x} = \mathbf{B}^{-1}(\overrightarrow{y})$ . Then the above expression becomes:  $\left(\frac{1}{\|\mathbf{A}^{-1}\|} - \|\mathbf{B} - \mathbf{A}\|\right) \|\mathbf{B}^{-1}(\overrightarrow{y})\| \leq \|\overrightarrow{y}\|$ . Because we are interested in  $\mathbf{B}$  close to  $\mathbf{A}$ , we can assume that  $\|\mathbf{B} - \mathbf{A}\| < \frac{1}{\|\mathbf{A}^{-1}\|}$ . Thus it is safe to divide by  $\frac{1}{\|\mathbf{A}^{-1}\|} - \|\mathbf{B} - \mathbf{A}\|$ . So, setting  $\overrightarrow{y} \in \mathbb{R}^n$  to be the vector satisfying  $\|\overrightarrow{y}\| \leq 1$  and  $\|\mathbf{B}^{-1}(\overrightarrow{y})\| = \|\mathbf{B}^{-1}\|$ , we have that:

$$\|\mathbf{B}^{-1}\| = \|\mathbf{B}^{-1}(\overrightarrow{y})\| \le \frac{\|\overrightarrow{y}\|}{\frac{1}{\|\mathbf{A}^{-1}\|} - \|\mathbf{B} - \mathbf{A}\|} \le \frac{1}{\frac{1}{\|\mathbf{A}^{-1}\|} - \|\mathbf{B} - \mathbf{A}\|} = \frac{\|\mathbf{A}^{-1}\|}{1 - \|\mathbf{A}^{-1}\| \|\mathbf{B} - \mathbf{A}\|}$$

 $\underline{\underline{\mathsf{Lemma}}}\text{: Given }\mathbf{A}\in L(Z,W)\text{, }\mathbf{B},\mathbf{C}\in L(Y,Z)\text{, and }\mathbf{D}\in L(X,Y)\text{,}$  we have that  $\mathbf{A}(\mathbf{B}+\mathbf{C})=\mathbf{AB}+\mathbf{AC}$  and  $(\mathbf{B}+\mathbf{C})\mathbf{D}=\mathbf{BD}+\mathbf{CD}$ .

Proof:

$$\circ \mathbf{A}((\mathbf{B} + \mathbf{C})(\overrightarrow{v})) = \mathbf{A}(\mathbf{B}(\overrightarrow{v}) + \mathbf{C}(\overrightarrow{v})) = \mathbf{A}(\mathbf{B}(\overrightarrow{v})) + \mathbf{A}(\mathbf{C}(\overrightarrow{v}))$$

$$\circ (\mathbf{B} + \mathbf{C})(\mathbf{D}(\overrightarrow{v})) = \mathbf{B}(\mathbf{D}(\overrightarrow{v})) + \mathbf{C}(\mathbf{D}(\overrightarrow{v}))$$

Based on the above lemma, we have that  ${\bf B}^{-1}-{\bf A}^{-1}={\bf B}^{-1}({\bf A}-{\bf B}){\bf A}^{-1}.$  So:

$$0 \le \|\mathbf{B}^{-1} - \mathbf{A}^{-1}\| = \|\mathbf{B}^{-1}(\mathbf{A} - \mathbf{B})\mathbf{A}^{-1}\|$$
  
$$\le \|\mathbf{B}^{-1}\|\|(\mathbf{A} - \mathbf{B})\|\|\mathbf{A}^{-1}\| \le \frac{\|\mathbf{A}^{-1}\|^2}{1 - \|\mathbf{A}^{-1}\|\|\mathbf{B} - \mathbf{A}\|}\|\mathbf{B} - \mathbf{A}\|$$

Finally, assume  $A \in \Omega'$ . This is fine because the mapping is automatically continuous at A if  $A \notin \Omega'$ . Then we have that:

$$\lim_{\mathbf{B}\to \mathbf{A}} \left( \frac{\|\mathbf{A}^{-1}\|^2}{1-\|\mathbf{A}^{-1}\|\|\mathbf{B}-\mathbf{A}\|} \|\mathbf{B}-\mathbf{A}\| \right) = \|\mathbf{A}^{-1}\|^2 \cdot 0 = 0.$$

So, 
$$0 \leq \lim_{\mathbf{B} \to \mathbf{A}} (\|\mathbf{B}^{-1} - \mathbf{A}^{-1}\|) \leq 0$$
.

This means that  $d(\mathbf{B}^{-1},\ \mathbf{A}^{-1})=\|\mathbf{B}^{-1}-\mathbf{A}^{-1}\|\to 0$  as  $\mathbf{B}\to\mathbf{A}$ . Or in other words:

$$\lim_{\mathbf{B}\to\mathbf{A}}(\mathbf{B}^{-1})=\mathbf{A}^{-1}.~\blacksquare$$

## Lecture 3: 4/9/2024

Let X and Y be vector spaces and fix two bases  $\{\vec{x}_1,\ldots,\vec{x}_n\}$  and  $\{\vec{y}_1,\ldots,\vec{y}_m\}$  of X and Y respectively. Then given any  $\mathbf{A}\in L(X,Y)$ , since  $\mathbf{A}(\vec{x}_j)\in Y$  for each  $j\in\{1,\ldots,n\}$ , we have that there are unique scalars  $a_{i,j}$  such that:

$$\mathbf{A}(\vec{x}_j) = \sum_{i=1}^m a_{i,j} \vec{y}_i$$

For convenience, we can visualize these numbers in an  $m \times n$  matrix:

$$[\mathbf{A}] = \begin{bmatrix} a_{1,1} & a_{1,2} & \cdots & a_{1,n} \\ a_{2,1} & a_{2,2} & \cdots & a_{2,n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{m,1} & a_{m,2} & \cdots & a_{m,n} \end{bmatrix}$$

Note that for each  $j \in \{1, ..., n\}$ , we have that the jth column of  $[\mathbf{A}]$  gives the coordinates of  $\mathbf{A}(\overrightarrow{x}_j)$  with respect to the basis  $\{\overrightarrow{y}_1, ..., \overrightarrow{y}_m\}$ . Thus, we call the vectors  $\mathbf{A}(\overrightarrow{x}_j)$  the <u>column vectors</u> of  $[\mathbf{A}]$ .

<u>Fact 1</u>: If  $\vec{x} \in X$ , then there are unique scalars  $c_1, \ldots, c_n$  such that  $\vec{x} = \sum_{j=1}^n c_j \vec{x}_j$ . Then, because **A** is linear:

$$\mathbf{A}(\overrightarrow{x}) = \mathbf{A} \left( \sum_{j=1}^{n} c_j \overrightarrow{x}_j \right) = \sum_{j=1}^{n} c_j \mathbf{A}(\overrightarrow{x}_j)$$
$$= \sum_{j=1}^{n} c_j \left( \sum_{i=1}^{m} a_{i,j} \overrightarrow{y}_i \right) = \sum_{i=1}^{m} \left( \sum_{j=1}^{n} c_j a_{i,j} \right) \overrightarrow{y}_i$$

Thus, the coordinates of  $\mathbf{A}(\overrightarrow{x})$  with respect to our basis of Y is given by the commonly defined matrix-vector product:

$$\begin{bmatrix} a_{1,1} & a_{1,2} & \cdots & a_{1,n} \\ a_{2,1} & a_{2,2} & \cdots & a_{2,n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{m,1} & a_{m,2} & \cdots & a_{m,n} \end{bmatrix} \begin{bmatrix} c_1 \\ c_2 \\ \vdots \\ c_n \end{bmatrix}$$

<u>Fact 2</u>: When we said how to generate an  $m \times n$  matrix  $[\mathbf{A}]$  for any  $A \in L(X,Y)$ , we were implicitely creating a mapping  $\phi: L(X,Y) \longrightarrow \mathcal{M}_{m \times n}(\mathbb{R})$  (the set of  $m \times n$  real matrices). Importantly, this map is one-to-one.

Let us define a mapping  $\varphi: \mathcal{M}_{m \times n}(\mathbb{R}) \longrightarrow L(X,Y)$  such that for any  $[\mathbf{B}] \in \mathcal{M}_{m \times n}(\mathbb{R})$  where

$$[\mathbf{B}] = \begin{bmatrix} b_{1,1} & b_{1,2} & \cdots & b_{1,n} \\ b_{2,1} & b_{2,2} & \cdots & b_{2,n} \\ \vdots & \vdots & \ddots & \vdots \\ b_{m,1} & b_{m,2} & \cdots & b_{m,n} \end{bmatrix}$$

...we define  $\varphi([\mathbf{B}]) \in L(X,Y)$  by  $\varphi([\mathbf{B}])(\overrightarrow{x}) = \sum_{i=1}^m \left(\sum_{j=1}^n c_j b_{i,j}\right) \overrightarrow{y}_i$  where  $c_1,\ldots,c_n$  are the coefficients such that  $\overrightarrow{x} = \sum_{j=1}^n c_j \overrightarrow{x}_j$ .

Firstly,  $\varphi([\mathbf{B}])$  is well defined because  $c_1, \ldots, c_n$  are unique with respect to our basis of X. Also  $\varphi([\mathbf{B}])$  is in fact linear because sums are linear.

Meanwhile, by fact 1 we know that for any  $\mathbf{A} \in L(X,Y), \ \mathbf{A} = \varphi(\phi(\mathbf{A})).$  At the same time, you can easily check that for any  $[\mathbf{B}] \in \mathcal{M}_{m \times n}(\mathbb{R})$ ,  $[\mathbf{B}] = \phi(\varphi([\mathbf{B}])).$  Hence,  $\varphi = \phi^{-1}.$ 

<u>Fact 3</u>: In addition to our bases for X and Y, fix  $\{\vec{z}_1,\ldots,\vec{z}_p\}$  as our basis for Z. Then given any  $\mathbf{B}\in L(Y,Z)$ , for each  $k\in 1,\ldots,m$  there are unique scalars  $b_{k,i}$  making up  $[\mathbf{B}]$  such that:

$$\mathbf{B}(\overrightarrow{y}_i) = \sum_{k=1}^p b_{k,i} \overrightarrow{z}_k$$

Therefore, for the linear map  $\mathbf{BA} \in L(X,Z)$ , we have that:

$$\mathbf{B}(\mathbf{A}(\overrightarrow{x_j})) = \mathbf{B}\left(\sum_{i=1}^m a_{i,j} \overrightarrow{y_i}\right) = \sum_{i=1}^m a_{i,j} \mathbf{B}(\overrightarrow{y_i})$$
$$= \sum_{i=1}^m a_{i,j} \left(\sum_{k=1}^p b_{k,i} \overrightarrow{z_k}\right) = \sum_{k=1}^p \left(\sum_{i=1}^m (a_{i,j} b_{k,i})\right) \overrightarrow{z_k}$$

So, the matrix  $[\mathbf{B}\mathbf{A}]$  generated by the linear map  $\mathbf{B}\mathbf{A} \in L(X,Z)$  is given by the commonly defined matrix-matrix product:  $[\mathbf{B}][\mathbf{A}]$ .

Now from a rigor point of view, we'd rather work with linear maps than matrices. This is because the defintion of a matrix depends on our choice of bases, whereas linear maps are defined independently of any bases. That said, matrices are too convenient to not be discussed.

One thing going forward though is that for any  $\mathbf{A} \in L(X,Y)$  and  $\overrightarrow{x} \in X$ , we shall abbreviate  $\mathbf{A}(\overrightarrow{x})$  as  $\mathbf{A}\overrightarrow{x}$  since it is already standard to do that for matrix-vector products.

Since an  $m \times n$  matrix can be thought of as a list of  $m \cdot n$  numbers, the "natural" norm to equip  $\mathcal{M}_{m \times n}(\mathbb{R})$  with is:

 $\|[\mathbf{A}]\|_F = \left(\sum_{i=1}^m \sum_{j=1}^n (a_{i,j})^2\right)^{\frac{1}{2}}$ 

Note on my notation:

Since I view  $|\cdot|$  as having already been reserved for the absolute value function, I am not going to use the same notation as Rudin and my professor use for this matrix norm. Rather, because this norm is also called the <u>Frobenius norm</u>, I shall denote it by  $||\cdot||_F$ .

Also, this is a valid norm for the same reasons that the vector Euclidean norm is a valid norm.

If we define  $d([\mathbf{B}], [\mathbf{A}]) = ||[\mathbf{B}] - [\mathbf{A}]||_F$ , then we can treat  $\mathcal{M}_{m \times n}(\mathbb{R})$  as a metric space with the metric d.

<u>Lemma 1</u>: Using the standard bases for  $\mathbb{R}^n$  and  $\mathbb{R}^m$ , we have that for any  $\mathbf{A} \in L(\mathbb{R}^n, \mathbb{R}^m)$  and its associated matrix  $[\mathbf{A}] \in \mathcal{M}_{m \times n}(\mathbb{R})$  with coefficients  $a_{i,j}$  for  $1 \le i \le m$  and  $1 \le j \le n$ :

$$\|\mathbf{A}\| \leq \|[\mathbf{A}]\|_F$$

Let  $\vec{x} = (x_1, \dots, x_n) \in \mathbb{R}^n$ . Then by the Cauchy-Schwarz inequality:

$$\|\mathbf{A}\,\vec{x}\|^2 = \sum_{i=1}^m \left(\sum_{j=1}^n a_{i,j} x_j\right)^2 \le \sum_{i=1}^m \left(\sum_{j=1}^n a_{i,j}^2 \cdot \sum_{j=1}^n x_j^2\right) = \|\vec{x}\|^2 \cdot \sum_{i=1}^m \sum_{j=1}^n a_{i,j}^2$$

So  $\|\mathbf{A}\overrightarrow{x}\|^2 \leq \|\overrightarrow{x}\|^2 \cdot \|[\mathbf{A}]\|_F^2$ . Or in other words,  $\|\mathbf{A}\|^2 \leq 1 \cdot \|[\mathbf{A}]\|_F^2$ .

<u>Lemma 2</u>: For any map  $\mathbf{A} \in L(X)$  and its associated matrix  $[\mathbf{A}] \in \mathcal{M}_{n \times n}(\mathbb{R})$ , we have that  $\mathbf{A}$  is invertible if and only if  $[\mathbf{A}]$  is invertible.

Proof:

( $\Longrightarrow$ ) If  ${\bf A}$  is invertible, then there exists  ${\bf A}^{-1}\in L(X)$  such that  ${\bf A}{\bf A}^{-1}\,\overrightarrow{x}=\overrightarrow{x}={\bf A}^{-1}{\bf A}\,\overrightarrow{x}$  for all  $\overrightarrow{x}\in X.$ 

<u>Proposition</u>: If  $\Omega$  is the set of invertible  $m\times n$  matrices, then  $\Omega$  is open in  $\mathcal{M}_{m\times n}(\mathbb{R})$ .

Proof:

Still using the standard bases for  $\mathbb{R}^n$  and  $\mathbb{R}^m$ , for any  $[\mathbf{A}] \in \mathcal{M}_{m \times n}(\mathbb{R})$ .

<u>Corollary</u>: If S is a matric space,  $a_{1,1},\ldots,a_{m,n}$  are real continuous functions on S, and for each  $p\in S$ ,  $\mathbf{A}_p$  is the linear map from  $\mathbb{R}^n$  to  $\mathbb{R}^m$  whose associated matrix is:

$$[\mathbf{A}_p] = \begin{bmatrix} a_{1,1} & a_{1,2} & \cdots & a_{1,n} \\ a_{2,1} & a_{2,2} & \cdots & a_{2,n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{m,1} & a_{m,2} & \cdots & a_{m,n} \end{bmatrix}$$

...then the mapping  $p\mapsto \mathbf{A}_p$  is a continuous mapping of S into  $L(\mathbb{R}^n,\mathbb{R}^m)$ .