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

Additionally, consider any  $\vec{x}_1, \vec{x}_2 \in X$ . Then 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, 
$$0 \le \lim_{\mathbf{B} \to \mathbf{A}} (\|\mathbf{B}^{-1} - \mathbf{A}^{-1}\|) \le 0$$
.

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

Suppose  $\{\vec{x}_1,\ldots,\vec{x}_n\}$  and  $\{\vec{y}_1,\ldots,\vec{y}_m\}$  are bases of the vector spaces X and Y respectively, and let  $\mathbf{A}\in L(X,Y)$ . Then for each  $j\in\{1,\ldots,n\}$ , since  $\mathbf{A}(\vec{x}_j)\in Y$ , there are unique coefficients  $a_{i,j}$  such that:

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

For convenience, we visualize these numbers in an  $\underline{m \times n}$  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  $\mathbf{A}(\vec{x}_j)$  with respect to the basis  $\{\vec{y}_1, ..., \vec{y}_m\}$ . Thus, we call the vectors  $\mathbf{A}(\vec{x}_j)$  the column vectors of  $[\mathbf{A}]$ .

<u>Fact</u>: There is a one-to-one correspondence between the set of  $m \times n$  real matrices and L(X,Y).

Firstly, let  $A \in L(X, Y)$ . Then we already saw above how to construct a matrix [A] from the linear mapping A.

Now observe if  $\overrightarrow{x} \in X$ , then  $\overrightarrow{x} = \sum_{j=1}^{n} c_j \overrightarrow{x_j}$ . Thus, because  $\mathbf{A}$  is linear:

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

Thus, we have an equation for  $\mathbf{A}(\overrightarrow{x})$  in terms of the components of  $[\mathbf{A}]$ . Needless to say, if we were instead starting out with an  $m \times n$  real matrix  $[\mathbf{B}] \in \mathcal{M}_{m \times n}(\mathbb{R})$  with components  $b_{i,j}$ , then we could define the linear map  $\mathbf{B} \in (L(X,Y))$  given by the rule:

$$\mathbf{B}(\overrightarrow{x}) = \sum_{i=1}^{m} \left( \sum_{j=1}^{n} c_j b_{i,j} \right) \overrightarrow{y}_i.$$