AN INTRODUCTION TO OPTIMIZATION

SOLUTIONS MANUAL

Fourth Edition

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A JOHN WILEY & SONS, INC., PUBLICATION

1. Methods of Proof and Some Notation

1.1 _____

А	В	not A	not B	A⇒B	$(\mathrm{not}\ B){\Rightarrow}(\mathrm{not}\ A)$
F	F	Т	Т	Т	Т
\mathbf{F}	Т	Т	\mathbf{F}	Т	Т
Т	\mathbf{F}	F	Т	F	\mathbf{F}
Т	Т	F	\mathbf{F}	Т	Т

1.2 _

A	В	not A	not B	A⇒B	not (A and (not B))
F	\mathbf{F}	Т	Т	Т	Т
\mathbf{F}	Т	Т	\mathbf{F}	Т	Т
Т	\mathbf{F}	F	Т	F	\mathbf{F}
Т	Т	F	F	Т	Т

1.3 _____

A	В	not (A and B)	not A	not B	(not A) or (not B))
F	F	Т	Т	Т	Т
\mathbf{F}	Т	Т	Т	\mathbf{F}	Т
Т	\mathbf{F}	Т	F	Т	Т
Т	Т	F	F	\mathbf{F}	\mathbf{F}

1.4 ____

A	В	A and B	A and (not B)	(A and B) or (A and (not B))
F	F	F	F	F
\mathbf{F}	Т	F	\mathbf{F}	\mathbf{F}
Т	F	F	Т	Т
Т	Т	Т	\mathbf{F}	Т

1.5 _____

The cards that you should turn over are 3 and A. The remaining cards are irrelevant to ascertaining the truth or falsity of the rule. The card with S is irrelevant because S is not a vowel. The card with 8 is not relevant because the rule does not say that if a card has an even number on one side, then it has a vowel on the other side.

Turning over the A card directly verifies the rule, while turning over the 3 card verifies the contraposition.

2. Vector Spaces and Matrices

2.1

We show this by contradiction. Suppose n < m. Then, the number of columns of A is n. Since rank A is the maximum number of linearly independent columns of A, then rank A cannot be greater than n < m, which contradicts the assumption that rank A = m.

2.2

 \Rightarrow : Since there exists a solution, then by Theorem 2.1, rank $\mathbf{A} = \operatorname{rank}[\mathbf{A}:\mathbf{b}]$. So, it remains to prove that rank $\mathbf{A} = n$. For this, suppose that rank $\mathbf{A} < n$ (note that it is impossible for rank $\mathbf{A} > n$ since \mathbf{A} has only n columns). Hence, there exists $\mathbf{y} \in \mathbb{R}^n$, $\mathbf{y} \neq \mathbf{0}$, such that $\mathbf{A}\mathbf{y} = \mathbf{0}$ (this is because the columns of

A are linearly dependent, and Ay is a linear combination of the columns of A). Let x be a solution to Ax = b. Then clearly $x + y \neq x$ is also a solution. This contradicts the uniqueness of the solution. Hence, rank A = n.

 \Leftarrow : By Theorem 2.1, a solution exists. It remains to prove that it is unique. For this, let x and y be solutions, i.e., Ax = b and Ay = b. Subtracting, we get A(x - y) = 0. Since rank A = n and A has n columns, then x - y = 0 and hence x = y, which shows that the solution is unique.

2.3_{-}

Consider the vectors $\bar{a}_i = [1, a_i^{\top}]^{\top} \in \mathbb{R}^{n+1}$, i = 1, ..., k. Since $k \ge n+2$, then the vectors $\bar{a}_1, ..., \bar{a}_k$ must be linearly independent in \mathbb{R}^{n+1} . Hence, there exist $\alpha_1, ..., \alpha_k$, not all zero, such that

$$\sum_{i=1}^k \alpha_i \boldsymbol{a}_i = \boldsymbol{0}.$$

The first component of the above vector equation is $\sum_{i=1}^{k} \alpha_i = 0$, while the last *n* components have the form $\sum_{i=1}^{k} \alpha_i a_i = 0$, completing the proof.

2.4

a. We first postmultiply M by the matrix

$$egin{bmatrix} oldsymbol{I}_k & oldsymbol{O} \ -oldsymbol{M}_{m-k,k} & oldsymbol{I}_{m-k} \end{bmatrix}$$

to obtain

$$egin{bmatrix} oldsymbol{M}_{m-k,k} & oldsymbol{I}_{m-k} \ oldsymbol{M}_{k,k} & oldsymbol{O} \end{bmatrix} egin{bmatrix} oldsymbol{I}_k & oldsymbol{O} \ -oldsymbol{M}_{m-k,k} & oldsymbol{I}_{m-k} \end{bmatrix} = egin{bmatrix} oldsymbol{O} & oldsymbol{I}_{m-k} \ oldsymbol{M}_{k,k} & oldsymbol{O} \end{bmatrix}$$

Note that the determinant of the postmultiplying matrix is 1. Next we postmultiply the resulting product by

$$egin{bmatrix} oldsymbol{O} & oldsymbol{I}_k \ oldsymbol{I}_{m-k} & oldsymbol{O} \end{bmatrix}$$

to obtain

$$egin{bmatrix} oldsymbol{O} & oldsymbol{I}_{m-k} \ oldsymbol{M}_{k,k} & oldsymbol{O} \end{bmatrix} egin{bmatrix} oldsymbol{O} & oldsymbol{I}_k \ oldsymbol{I}_{m-k} & oldsymbol{O} \end{bmatrix} = egin{bmatrix} oldsymbol{I}_k & oldsymbol{O} \ oldsymbol{O} & oldsymbol{M}_{k,k} \end{bmatrix}.$$

Notice that

where

$$\begin{split} \det \boldsymbol{M} &= \det \left(\begin{bmatrix} \boldsymbol{I}_k & \boldsymbol{O} \\ \boldsymbol{O} & \boldsymbol{M}_{k,k} \end{bmatrix} \right) \det \left(\begin{bmatrix} \boldsymbol{O} & \boldsymbol{I}_k \\ \boldsymbol{I}_{m-k} & \boldsymbol{O} \end{bmatrix} \right), \\ & \det \left(\begin{bmatrix} \boldsymbol{O} & \boldsymbol{I}_k \\ \boldsymbol{I}_{m-k} & \boldsymbol{O} \end{bmatrix} \right) = \pm 1. \end{split}$$

The above easily follows from the fact that the determinant changes its sign if we interchange columns, as discussed in Section 2.2. Moreover,

$$\det \left(\begin{bmatrix} \boldsymbol{I}_k & \boldsymbol{O} \\ \boldsymbol{O} & \boldsymbol{M}_{k,k} \end{bmatrix} \right) = \det(\boldsymbol{I}_k) \det(\boldsymbol{M}_{k,k}) = \det(\boldsymbol{M}_{k,k}).$$

Hence,

$$\det \boldsymbol{M} = \pm \det \boldsymbol{M}_{k,k}.$$

b. We can see this on the following examples. We assume, without loss of generality that $M_{m-k,k} = O$ and let $M_{k,k} = 2$. Thus k = 1. First consider the case when m = 2. Then we have

$$\boldsymbol{M} = \begin{bmatrix} \boldsymbol{O} & \boldsymbol{I}_{m-k} \\ \boldsymbol{M}_{k,k} & \boldsymbol{O} \end{bmatrix} = \begin{bmatrix} 0 & 1 \\ 2 & 0 \end{bmatrix}.$$

Thus,

$$\det \boldsymbol{M} = -2 = \det \left(-\boldsymbol{M}_{k,k} \right)$$

Next consider the case when m = 3. Then

$$\det \begin{bmatrix} \boldsymbol{O} & \boldsymbol{I}_{m-k} \\ \boldsymbol{M}_{k,k} & \boldsymbol{O} \end{bmatrix} = \det \begin{bmatrix} 0 & \vdots & 1 & 0 \\ 0 & \vdots & 0 & 1 \\ \dots & \dots & \dots & \dots \\ 2 & \vdots & 0 & 0 \end{bmatrix} = 2 \neq \det \left(-\boldsymbol{M}_{k,k} \right).$$

Therefore, in general,

 $\det \boldsymbol{M} \neq \det \left(-\boldsymbol{M}_{k,k}\right)$

However, when k = m/2, that is, when all sub-matrices are square and of the same dimension, then it is true that

$$\det \boldsymbol{M} = \det \left(-\boldsymbol{M}_{k,k}\right).$$

See [121].

 $\mathbf{2.5}$ _

Let

$$M = egin{bmatrix} A & B \ C & D \end{bmatrix}$$

and suppose that each block is $k \times k$. John R. Silvester [121] showed that if at least one of the blocks is equal to O (zero matrix), then the desired formula holds. Indeed, if a row or column block is zero, then the determinant is equal to zero as follows from the determinant's properties discussed Section 2.2. That is, if A = B = O, or A = C = O, and so on, then obviously det M = 0. This includes the case when any three or all four block matrices are zero matrices.

If B = O or C = O then

$$\det M = \det \begin{bmatrix} A & B \\ C & D \end{bmatrix} = \det (AD).$$

The only case left to analyze is when A = O or D = O. We will show that in either case,

 $\det \boldsymbol{M} = \det \left(-\boldsymbol{B}\boldsymbol{C} \right).$

Without loss of generality suppose that D = O. Following arguments of John R. Silvester [121], we premultiply M by the product of three matrices whose determinants are unity:

$$\begin{bmatrix} I_k & -I_k \\ O & I_k \end{bmatrix} \begin{bmatrix} I_k & O \\ I_k & I_k \end{bmatrix} \begin{bmatrix} I_k & -I_k \\ O & I_k \end{bmatrix} \begin{bmatrix} A & B \\ C & O \end{bmatrix} = \begin{bmatrix} -C & O \\ A & B \end{bmatrix}$$

Hence,

$$\det \begin{bmatrix} A & B \\ C & O \end{bmatrix} = \begin{bmatrix} -C & O \\ A & B \end{bmatrix}$$
$$= \det (-C) \det B$$
$$= \det (-I_k) \det C \det B.$$

Thus we have

$$\det \begin{bmatrix} A & B \\ C & O \end{bmatrix} = \det (-BC) = \det (-CB)$$

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We represent the given system of equations in the form Ax = b, where

$$m{A} = egin{bmatrix} 1 & 1 & 2 & 1 \ 1 & -2 & 0 & -1 \end{bmatrix}, \quad m{x} = egin{bmatrix} x_1 \ x_2 \ x_3 \ x_4 \end{bmatrix}, \quad ext{and} \quad m{b} = egin{bmatrix} 1 \ -2 \end{bmatrix}.$$

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Using elementary row operations yields

2.6 _____

$$\boldsymbol{A} = \begin{bmatrix} 1 & 1 & 2 & 1 \\ 1 & -2 & 0 & -1 \end{bmatrix} \rightarrow \begin{bmatrix} 1 & 1 & 2 & 1 \\ 0 & -3 & -2 & -2 \end{bmatrix}, \text{ and}$$
$$[\boldsymbol{A}, \boldsymbol{b}] = \begin{bmatrix} 1 & 1 & 2 & 1 & 1 \\ 1 & -2 & 0 & -1 & -2 \end{bmatrix} \rightarrow \begin{bmatrix} 1 & 1 & 2 & 1 & 1 \\ 0 & -3 & -2 & -2 & -3 \end{bmatrix},$$

from which rank A = 2 and rank [A, b] = 2. Therefore, by Theorem 2.1, the system has a solution.

We next represent the system of equations as

$$\begin{bmatrix} 1 & 1 \\ 1 & -2 \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \end{bmatrix} = \begin{bmatrix} 1 - 2x_3 - x_4 \\ -2 + x_4 \end{bmatrix}$$

Assigning arbitrary values to x_3 and x_4 ($x_3 = d_3$, $x_4 = d_4$), we get

$$\begin{bmatrix} x_1 \\ x_2 \end{bmatrix} = \begin{bmatrix} 1 & 1 \\ 1 & -2 \end{bmatrix}^{-1} \begin{bmatrix} 1 - 2x_3 - x_4 \\ -2 + x_4 \end{bmatrix}$$
$$= -\frac{1}{3} \begin{bmatrix} -2 & -1 \\ -1 & 1 \end{bmatrix} \begin{bmatrix} 1 - 2x_3 - x_4 \\ -2 + x_4 \end{bmatrix}$$
$$= \begin{bmatrix} -\frac{4}{3}d_3 - \frac{1}{3}d_4 \\ 1 - \frac{2}{3}d_3 - \frac{2}{3}d_4 \end{bmatrix}.$$

Therefore, a general solution is

$$\begin{bmatrix} x_1 \\ x_2 \\ x_3 \\ x_4 \end{bmatrix} = \begin{bmatrix} -\frac{4}{3}d_3 - \frac{1}{3}d_4 \\ 1 - \frac{2}{3}d_3 - \frac{2}{3}d_4 \\ d_3 \\ d_4 \end{bmatrix} = \begin{bmatrix} -\frac{4}{3} \\ -\frac{2}{3} \\ 1 \\ 0 \end{bmatrix} d_3 + \begin{bmatrix} -\frac{1}{3} \\ -\frac{2}{3} \\ 0 \\ 1 \end{bmatrix} d_4 + \begin{bmatrix} 0 \\ 1 \\ 0 \\ 0 \end{bmatrix},$$

where d_3 and d_4 are arbitrary values.

2.7 ____

1. Apply the definition of |-a|:

$$|-a| = \begin{cases} -a & \text{if } -a > 0\\ 0 & \text{if } -a = 0\\ -(-a) & \text{if } -a < 0 \end{cases}$$
$$= \begin{cases} -a & \text{if } a < 0\\ 0 & \text{if } a = 0\\ a & \text{if } a > 0 \end{cases}$$
$$= |a|.$$

2. If $a \ge 0$, then |a| = a. If a < 0, then |a| = -a > 0 > a. Hence $|a| \ge a$. On the other hand, $|-a| \ge -a$ (by the above). Hence, $a \ge -|-a| = -|a|$ (by property 1).

3. We have four cases to consider. First, if $a, b \ge 0$, then $a + b \ge 0$. Hence, |a + b| = a + b = |a| + |b|. Second, if $a, b \ge 0$, then $a + b \le 0$. Hence |a + b| = -(a + b) = -a - b = |a| + |b|. Third, if $a \ge 0$ and $b \le 0$, then we have two further subcases:

1. If $a + b \ge 0$, then $|a + b| = a + b \le |a| + |b|$.

2. If $a + b \le 0$, then $|a + b| = -a - b \le |a| + |b|$.

The fourth case, $a \le 0$ and $b \ge 0$, is identical to the third case, with a and b interchanged. 4. We first show $|a - b| \le |a| + |b|$. We have

$$\begin{aligned} |a-b| &= |a+(-b)| \\ &\leq |a|+|-b| \quad \text{by property 3} \\ &= |a|+|b| \quad \text{by property 1.} \end{aligned}$$

To show $||a| - |b|| \le |a - b|$, we note that $|a| = |a - b + b| \le |a - b| + |b|$, which implies $|a| - |b| \le |a - b|$. On the other hand, from the above we have $|b| - |a| \le |b - a| = |a - b|$ by property 1. Therefore, $||a| - |b|| \le |a - b|$.

5. We have four cases. First, if $a, b \ge 0$, we have $ab \ge 0$ and hence |ab| = ab = |a||b|. Second, if $a, b \le 0$, we have $ab \ge 0$ and hence |ab| = ab = (-a)(-b) = |a||b|. Third, if $a \le 0$, $b \le 0$, we have $ab \le 0$ and hence |ab| = -ab = a(-b) = |a||b|. The fourth case, $a \le 0$ and $b \ge 0$, is identical to the third case, with a and b interchanged.

6. We have

$$|a+b| \leq |a|+|b|$$
 by property 3
 $< c+d.$

7. \Rightarrow : By property 2, $-a \leq |a|$ and $a \leq |a$. Therefore, |a| < b implies $-a \leq |a| < b$ and $a \leq |a| < b$. \Leftarrow : If $a \geq 0$, then |a| = a < b. If a < 0, then |a| = -a < b.

For the case when "<" is replaced by " \leq ", we simply repeat the above proof with "<" replaced by " \leq ". 8. This is simply the negation of property 7 (apply DeMorgan's Law).

2.8 -

Observe that we can represent $\langle \boldsymbol{x}, \boldsymbol{y} \rangle_2$ as

$$\langle \boldsymbol{x}, \boldsymbol{y} \rangle_2 = \boldsymbol{x}^\top \begin{bmatrix} 2 & 3 \\ 3 & 5 \end{bmatrix} \boldsymbol{y} = (\boldsymbol{Q} \boldsymbol{x})^\top (\boldsymbol{Q} \boldsymbol{y}) = \boldsymbol{x}^\top \boldsymbol{Q}^2 \boldsymbol{y},$$

where

$$oldsymbol{Q} = egin{bmatrix} 1 & 1 \ 1 & 2 \end{bmatrix}$$

Note that the matrix $\boldsymbol{Q} = \boldsymbol{Q}^{\top}$ is nonsingular.

1. Now, $\langle \boldsymbol{x}, \boldsymbol{x} \rangle_2 = (\boldsymbol{Q}\boldsymbol{x})^\top (\boldsymbol{Q}\boldsymbol{x}) = \|\boldsymbol{Q}\boldsymbol{x}\|^2 \ge 0$, and

$$\langle \boldsymbol{x}, \boldsymbol{x} \rangle_2 = 0 \quad \Leftrightarrow \quad \| \boldsymbol{Q} \boldsymbol{x} \|^2 = 0$$

 $\Leftrightarrow \quad \boldsymbol{Q} \boldsymbol{x} = \mathbf{0}$
 $\Leftrightarrow \quad \boldsymbol{x} = \mathbf{0}$

since Q is nonsingular.

2. $\langle \boldsymbol{x}, \boldsymbol{y} \rangle_2 = (\boldsymbol{Q} \boldsymbol{x})^\top (\boldsymbol{Q} \boldsymbol{y}) = (\boldsymbol{Q} \boldsymbol{y})^\top (\boldsymbol{Q} \boldsymbol{x}) = \langle \boldsymbol{y}, \boldsymbol{x} \rangle_2.$

3. We have

$$egin{array}{rcl} \langle m{x}+m{y},m{z}
angle_2&=&(m{x}+m{y})^{ op}m{Q}^2m{z}\ &=&m{x}^{ op}m{Q}^2m{z}+m{y}^{ op}m{Q}^2m{z}\ &=&\langlem{x},m{z}
angle_2+m{y}^{ op}m{Q}^2m{z}. \end{array}$$

4.
$$\langle r\boldsymbol{x}, \boldsymbol{y} \rangle_2 = (r\boldsymbol{x})^\top \boldsymbol{Q}^2 \boldsymbol{y} = r\boldsymbol{x}^\top \boldsymbol{Q}^2 \boldsymbol{y} = r \langle \boldsymbol{x}, \boldsymbol{y} \rangle_2.$$

2.9_{-}

We have $\|\boldsymbol{x}\| = \|(\boldsymbol{x} - \boldsymbol{y}) + \boldsymbol{y}\| \le \|\boldsymbol{x} - \boldsymbol{y}\| + \|\boldsymbol{y}\|$ by the Triangle Inequality. Hence, $\|\boldsymbol{x}\| - \|\boldsymbol{y}\| \le \|\boldsymbol{x} - \boldsymbol{y}\|$. On the other hand, from the above we have $\|\boldsymbol{y}\| - \|\boldsymbol{x}\| \le \|\boldsymbol{y} - \boldsymbol{x}\| = \|\boldsymbol{x} - \boldsymbol{y}\|$. Combining the two inequalities, we obtain $\|\|\boldsymbol{x}\| - \|\boldsymbol{y}\| \le \|\boldsymbol{x} - \boldsymbol{y}\|$.

2.10 –

Let $\epsilon > 0$ be given. Set $\delta = \epsilon$. Hence, if $||\mathbf{x} - \mathbf{y}|| < \delta$, then by Exercise 2.9, $|||\mathbf{x}|| - ||\mathbf{y}|| \le ||\mathbf{x} - \mathbf{y}|| < \delta = \epsilon$.

3. Transformations

3.1 _

Let v be the vector such that x are the coordinates of v with respect to $\{e_1, e_2, \ldots, e_n\}$, and x' are the coordinates of v with respect to $\{e'_1, e'_2, \ldots, e'_n\}$. Then,

$$\boldsymbol{v} = x_1 \boldsymbol{e}_1 + \cdots + x_n \boldsymbol{e}_n = [\boldsymbol{e}_1, \dots, \boldsymbol{e}_n] \boldsymbol{x},$$

and

$$\boldsymbol{v} = x_1' \boldsymbol{e}_1' + \dots + x_n' \boldsymbol{e}_n' = [\boldsymbol{e}_1', \dots, \boldsymbol{e}_n'] \boldsymbol{x}'$$

Hence,

$$[oldsymbol{e}_1,\ldots,oldsymbol{e}_n]oldsymbol{x}=[oldsymbol{e}_1',\ldots,oldsymbol{e}_n']oldsymbol{x}'$$

which implies

$$oldsymbol{x}' = [oldsymbol{e}_1',\ldots,oldsymbol{e}_n']^{-1}[oldsymbol{e}_1,\ldots,oldsymbol{e}_n]oldsymbol{x} = oldsymbol{T}oldsymbol{x}.$$

3.2 _

a. We have

$$[\boldsymbol{e}_1', \boldsymbol{e}_2', \boldsymbol{e}_3'] = [\boldsymbol{e}_1, \boldsymbol{e}_2, \boldsymbol{e}_3] \begin{bmatrix} 1 & 2 & 4 \\ 3 & -1 & 5 \\ -4 & 5 & 3 \end{bmatrix}.$$

Therefore,

$$\boldsymbol{T} = [\boldsymbol{e}_1', \boldsymbol{e}_2', \boldsymbol{e}_3']^{-1}[\boldsymbol{e}_1, \boldsymbol{e}_2, \boldsymbol{e}_3] = \begin{bmatrix} 1 & 2 & 4 \\ 3 & -1 & 5 \\ -4 & 5 & 3 \end{bmatrix}^{-1} = \frac{1}{42} \begin{bmatrix} 28 & -14 & -14 \\ 29 & -19 & -7 \\ -11 & 13 & 7 \end{bmatrix}.$$

b. We have

$$[\boldsymbol{e}_1, \boldsymbol{e}_2, \boldsymbol{e}_3] = [\boldsymbol{e}_1', \boldsymbol{e}_2', \boldsymbol{e}_3'] \begin{bmatrix} 1 & 2 & 3 \\ 1 & -1 & 0 \\ 3 & 4 & 5 \end{bmatrix}.$$

Therefore,

$$\boldsymbol{T} = \begin{bmatrix} 1 & 2 & 3 \\ 1 & -1 & 0 \\ 3 & 4 & 5 \end{bmatrix}$$

3.3 _

We have

$$[\boldsymbol{e}_1, \boldsymbol{e}_2, \boldsymbol{e}_3] = [\boldsymbol{e}_1', \boldsymbol{e}_2', \boldsymbol{e}_3'] \begin{bmatrix} 2 & 2 & 3 \\ 1 & -1 & 0 \\ -1 & 2 & 1 \end{bmatrix}.$$

Therefore, the transformation matrix from $\{e'_1, e'_2, e'_3\}$ to $\{e_1, e_2, e_3\}$ is

$$\boldsymbol{T} = \begin{bmatrix} 2 & 2 & 3 \\ 1 & -1 & 0 \\ -1 & 2 & 1 \end{bmatrix},$$

Now, consider a linear transformation $L : \mathbb{R}^3 \to \mathbb{R}^3$, and let A be its representation with respect to $\{e_1, e_2, e_3\}$, and B its representation with respect to $\{e_1', e_2', e_3'\}$. Let y = Ax and y' = Bx'. Then,

$$y' = Ty = T(Ax) = TA(T^{-1}x') = (TAT^{-1})x'.$$

Hence, the representation of the linear transformation with respect to $\{ \bm{e}_1', \bm{e}_2', \bm{e}_3' \}$ is

$$\boldsymbol{B} = \boldsymbol{T}\boldsymbol{A}\boldsymbol{T}^{-1} = \begin{bmatrix} 3 & -10 & -8 \\ -1 & 8 & 4 \\ 2 & -13 & -7 \end{bmatrix}.$$

3.4 ____

We have

$$[\boldsymbol{e}_1', \boldsymbol{e}_2', \boldsymbol{e}_3', \boldsymbol{e}_4'] = [\boldsymbol{e}_1, \boldsymbol{e}_2, \boldsymbol{e}_3, \boldsymbol{e}_4] \begin{bmatrix} 1 & 1 & 1 & 1 \\ 0 & 1 & 1 & 1 \\ 0 & 0 & 1 & 1 \\ 0 & 0 & 0 & 1 \end{bmatrix}.$$

Therefore, the transformation matrix from $\{e_1, e_2, e_3, e_4\}$ to $\{e'_1, e'_2, e'_3, e'_4\}$ is

$$\boldsymbol{T} = \begin{bmatrix} 1 & 1 & 1 & 1 \\ 0 & 1 & 1 & 1 \\ 0 & 0 & 1 & 1 \\ 0 & 0 & 0 & 1 \end{bmatrix}^{-1} = \begin{bmatrix} 1 & -1 & 0 & 0 \\ 0 & 1 & -1 & 0 \\ 0 & 0 & 1 & -1 \\ 0 & 0 & 0 & 1 \end{bmatrix}.$$

Now, consider a linear transformation $L: \mathbb{R}^4 \to \mathbb{R}^4$, and let **A** be its representation with respect to $\{e_1, e_2, e_3, e_4\}$, and B its representation with respect to $\{e'_1, e'_2, e'_3, e'_4\}$. Let y = Ax and y' = Bx'. Then,

$$y' = Ty = T(Ax) = TA(T^{-1}x') = (TAT^{-1})x'.$$

Therefore,

$$\boldsymbol{B} = \boldsymbol{T}\boldsymbol{A}\boldsymbol{T}^{-1} = \begin{bmatrix} 5 & 3 & 4 & 3 \\ -3 & -2 & -1 & -2 \\ -1 & 0 & -1 & -2 \\ 1 & 1 & 1 & 4 \end{bmatrix}.$$

 $\begin{array}{c} \textbf{3.5} \\ \textbf{Let } \{ \boldsymbol{v}_1, \boldsymbol{v}_2, \boldsymbol{v}_3, \boldsymbol{v}_4 \} \text{ be a set of linearly independent eigenvectors of } \boldsymbol{A} \text{ corresponding to the eigenvalues } \lambda_1, \end{array}$ λ_2 , λ_3 , and λ_4 . Let $\boldsymbol{T} = [\boldsymbol{v}_1, \boldsymbol{v}_2, \boldsymbol{v}_3, \boldsymbol{v}_4]$. Then,

$$\begin{aligned} \boldsymbol{AT} &= \boldsymbol{A}[\boldsymbol{v}_1, \boldsymbol{v}_2, \boldsymbol{v}_3, \boldsymbol{v}_4] = [\boldsymbol{Av}_1, \boldsymbol{Av}_2, \boldsymbol{Av}_3, \boldsymbol{Av}_4] \\ &= [\lambda_1 \boldsymbol{v}_1, \lambda_2 \boldsymbol{v}_2, \lambda_3 \boldsymbol{v}_3, \lambda_4 \boldsymbol{v}_4] = [\boldsymbol{v}_1, \boldsymbol{v}_2, \boldsymbol{v}_3, \boldsymbol{v}_4] \begin{bmatrix} \lambda_1 & 0 & 0 & 0 \\ 0 & \lambda_2 & 0 & 0 \\ 0 & 0 & \lambda_3 & 0 \\ 0 & 0 & 0 & \lambda_4 \end{bmatrix} \end{aligned}$$

Hence,

$$oldsymbol{AT} = oldsymbol{T} egin{bmatrix} \lambda_1 & 0 & 0 \ 0 & \lambda_2 & 0 \ 0 & 0 & \lambda_3 \end{bmatrix},$$

or

$$\boldsymbol{T}^{-1}\boldsymbol{A}\boldsymbol{T} = \begin{bmatrix} \lambda_1 & 0 & 0 \\ 0 & \lambda_2 & 0 \\ 0 & 0 & \lambda_3 \end{bmatrix}.$$

Therefore, the linear transformation has a diagonal matrix form with respect to the basis formed by a linearly independent set of eigenvectors.

Because

$$\det(\mathbf{A}) = (\lambda - 2)(\lambda - 3)(\lambda - 1)(\lambda + 1),$$

the eigenvalues are $\lambda_1 = 2$, $\lambda_2 = 3$, $\lambda_3 = 1$, and $\lambda_4 = -1$.

From $Av_i = \lambda_i v_i$, where $v_i \neq 0$ (i = 1, 2, 3), the corresponding eigenvectors are

$$\boldsymbol{v}_1 = \begin{bmatrix} 0\\0\\1\\0 \end{bmatrix}, \quad \boldsymbol{v}_2 = \begin{bmatrix} 0\\0\\1\\1 \end{bmatrix}, \quad \boldsymbol{v}_3 = \begin{bmatrix} 0\\2\\-9\\1 \end{bmatrix}, \text{and} \quad \boldsymbol{v}_4 = \begin{bmatrix} 24\\-12\\1\\9 \end{bmatrix}.$$

Therefore, the basis we are interested in is

$$\{\boldsymbol{v}_1, \boldsymbol{v}_2, \boldsymbol{v}_3\} = \left\{ \begin{bmatrix} 0\\0\\1\\1 \end{bmatrix}, \begin{bmatrix} 0\\0\\1\\1 \end{bmatrix}, \begin{bmatrix} 0\\0\\1\\1 \end{bmatrix}, \begin{bmatrix} 0\\2\\-9\\1 \end{bmatrix}, \begin{bmatrix} 24\\-12\\1\\9 \end{bmatrix} \right\}$$

3.6 ____

Suppose v_1, \ldots, v_n are eigenvectors of A corresponding to $\lambda_1, \ldots, \lambda_n$, respectively. Then, for each i = $1, \ldots, n$, we have

$$(\boldsymbol{I}_n - \boldsymbol{A})\boldsymbol{v}_i = \boldsymbol{v}_i - \boldsymbol{A}\boldsymbol{v}_i = \boldsymbol{v}_i - \lambda_i \boldsymbol{v}_i = (1 - \lambda_i)\boldsymbol{v}_i$$

which shows that $1 - \lambda_1, \ldots, 1 - \lambda_n$ are the eigenvalues of $I_n - A$.

Alternatively, we may write the characteristic polynomial of $\boldsymbol{I}_n-\boldsymbol{A}$ as

$$\pi_{\boldsymbol{I}_n-\boldsymbol{A}}(1-\lambda) = \det((1-\lambda)\boldsymbol{I}_n - (\boldsymbol{I}_n - \boldsymbol{A})) = \det(-[\lambda \boldsymbol{I}_n - \boldsymbol{A}]) = (-1)^n \pi_{\boldsymbol{A}}(\lambda),$$

which shows the desired result.

3.7 Let $x, y \in \mathcal{V}^{\perp}$, and $\alpha, \beta \in \mathbb{R}$. To show that \mathcal{V}^{\perp} is a subspace, we need to show that $\alpha x + \beta y \in \mathcal{V}^{\perp}$. For this, let v be any vector in \mathcal{V} . Then,

$$\boldsymbol{v}^{\top}(\alpha \boldsymbol{x} + \beta \boldsymbol{y}) = \alpha \boldsymbol{v}^{\top} \boldsymbol{x} + \beta \boldsymbol{v}^{\top} \boldsymbol{y} = 0,$$

since $\boldsymbol{v}^{\top}\boldsymbol{x} = \boldsymbol{v}^{\top}\boldsymbol{y} = 0$ by definition.

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3.8 _

The null space of A is $\mathcal{N}(A) = \{x \in \mathbb{R}^3 : Ax = 0\}$. Using elementary row operations and back-substitution, we can solve the system of equations:

$$\begin{bmatrix} 4 & -2 & 0 \\ 2 & 1 & -1 \\ 2 & -3 & 1 \end{bmatrix} \rightarrow \begin{bmatrix} 4 & -2 & 0 \\ 0 & 2 & -1 \\ 0 & -2 & 1 \end{bmatrix} \rightarrow \begin{bmatrix} 4 & -2 & 0 \\ 0 & 2 & -1 \\ 0 & 0 & 0 \end{bmatrix} \qquad \Rightarrow \qquad \begin{aligned} 4x_1 - 2x_2 &= 0 \\ 2x_2 - x_3 &= 0 \end{aligned}$$
$$\Rightarrow \qquad x_2 = \frac{1}{2}x_3, \qquad x_1 = \frac{1}{2}x_2 = \frac{1}{4}x_3 \qquad \Rightarrow \qquad x = \begin{bmatrix} x_1 \\ x_2 \\ x_3 \end{bmatrix} = \begin{bmatrix} \frac{1}{4} \\ \frac{1}{2} \\ 1 \end{bmatrix} x_3.$$

Therefore,

$$\mathcal{N}(\boldsymbol{A}) = \left\{ \begin{bmatrix} 1\\ 2\\ 4 \end{bmatrix} c : c \in \mathbb{R} \right\}.$$

3.9 _

Let $x, y \in \mathcal{R}(A)$, and $\alpha, \beta \in \mathbb{R}$. Then, there exists v, u such that x = Av and y = Au. Thus,

$$\alpha \boldsymbol{x} + \beta \boldsymbol{y} = \alpha \boldsymbol{A} \boldsymbol{v} + \beta \boldsymbol{A} \boldsymbol{u} = \boldsymbol{A} (\alpha \boldsymbol{v} + \beta \boldsymbol{u}).$$

Hence, $\alpha \boldsymbol{x} + \beta \boldsymbol{y} \in \mathcal{R}(\boldsymbol{A})$, which shows that $\mathcal{R}(\boldsymbol{A})$ is a subspace.

Let $x, y \in \mathcal{N}(A)$, and $\alpha, \beta \in \mathbb{R}$. Then, Ax = 0 and Ay = 0. Thus,

$$\boldsymbol{A}(\alpha \boldsymbol{x} + \beta \boldsymbol{y}) = \alpha \boldsymbol{A} \boldsymbol{x} + \beta \boldsymbol{A} \boldsymbol{y} = \boldsymbol{0}$$

Hence, $\alpha \boldsymbol{x} + \beta \boldsymbol{y} \in \mathcal{N}(\boldsymbol{A})$, which shows that $\mathcal{N}(\boldsymbol{A})$ is a subspace.

3.10_{-}

Let $\boldsymbol{v} \in \mathcal{R}(\boldsymbol{B})$, i.e., $\boldsymbol{v} = \boldsymbol{B}\boldsymbol{x}$ for some \boldsymbol{x} . Consider the matrix $[\boldsymbol{A} \ \boldsymbol{v}]$. Then, $\mathcal{N}(\boldsymbol{A}^{\top}) = \mathcal{N}([\boldsymbol{A} \ \boldsymbol{v}]^{\top})$, since if $\boldsymbol{u} \in \mathcal{N}(\boldsymbol{A}^{\top})$, then $\boldsymbol{u} \in \mathcal{N}(\boldsymbol{B}^{\top})$ by assumption, and hence $\boldsymbol{u}^{\top}\boldsymbol{v} = \boldsymbol{u}^{\top}\boldsymbol{B}\boldsymbol{x} = \boldsymbol{x}^{\top}\boldsymbol{B}^{\top}\boldsymbol{u} = \boldsymbol{0}$. Now,

$$\dim \mathcal{R}(\mathbf{A}) + \dim \mathcal{N}(\mathbf{A}^{\top}) = m$$

and

$$\dim \mathcal{R}([\boldsymbol{A} \ \boldsymbol{v}]) + \dim \mathcal{N}([\boldsymbol{A} \ \boldsymbol{v}]^{\top}) = m.$$

Since dim $\mathcal{N}(\mathbf{A}^{\top}) = \dim \mathcal{N}([\mathbf{A} \ \mathbf{v}]^{\top})$, then we have dim $\mathcal{R}(\mathbf{A}) = \dim \mathcal{R}([\mathbf{A} \ \mathbf{v}])$. Hence, \mathbf{v} is a linear combination of the columns of \mathbf{A} , i.e., $\mathbf{v} \in \mathcal{R}(\mathbf{A})$, which completes the proof.

3.11

We first show $V \subset (V^{\perp})^{\perp}$. Let $v \in V$, and u any element of V^{\perp} . Then $u^{\top}v = v^{\top}u = 0$. Therefore, $v \in (V^{\perp})^{\perp}$.

We now show $(\mathbf{V}^{\perp})^{\perp} \subset \mathbf{V}$. Let $\{a_1, \ldots, a_k\}$ be a basis for \mathbf{V} , and $\{b_1, \ldots, b_l\}$ a basis for $(\mathbf{V}^{\perp})^{\perp}$. Define $\mathbf{A} = [a_1 \cdots a_k]$ and $\mathbf{B} = [b_1 \cdots b_l]$, so that $\mathbf{V} = \mathcal{R}(\mathbf{A})$ and $(\mathbf{V}^{\perp})^{\perp} = \mathcal{R}(\mathbf{B})$. Hence, it remains to show that $\mathcal{R}(\mathbf{B}) \subset \mathcal{R}(\mathbf{A})$. Using the result of Exercise 3.10, it suffices to show that $\mathcal{N}(\mathbf{A}^{\top}) \subset \mathcal{N}(\mathbf{B}^{\top})$. So let $\mathbf{x} \in \mathcal{N}(\mathbf{A}^{\top})$, which implies that $\mathbf{x} \in \mathcal{R}(\mathbf{A})^{\perp} = \mathbf{V}^{\perp}$, since $\mathcal{R}(\mathbf{A})^{\perp} = \mathcal{N}(\mathbf{A}^{\top})$. Hence, for all \mathbf{y} , we have $(\mathbf{B}\mathbf{y})^{\top}\mathbf{x} = \mathbf{0} = \mathbf{y}^{\top}\mathbf{B}^{\top}\mathbf{x}$, which implies that $\mathbf{B}^{\top}\mathbf{x} = \mathbf{0}$. Therefore, $\mathbf{x} \in \mathcal{N}(\mathbf{B}^{\top})$, which completes the proof. 3.12

Let $\boldsymbol{w} \in \mathcal{W}^{\perp}$, and \boldsymbol{y} be any element of \mathcal{V} . Since $\mathcal{V} \subset \mathcal{W}$, then $\boldsymbol{y} \in \mathcal{W}$. Therefore, by definition of \boldsymbol{w} , we have $\boldsymbol{w}^{\top}\boldsymbol{y} = 0$. Therefore, $\boldsymbol{w} \in \mathcal{V}^{\perp}$.

3.13 -

Let $r = \dim \mathcal{V}$. Let v_1, \ldots, v_r be a basis for \mathcal{V} , and V the matrix whose *i*th column is v_i . Then, clearly $\mathcal{V} = \mathcal{R}(V)$.

Let $\mathbf{u}_1, \ldots, \mathbf{u}_{n-r}$ be a basis for \mathcal{V}^{\perp} , and \mathbf{U} the matrix whose *i*th row is \mathbf{u}_i^{\top} . Then, $\mathcal{V}^{\perp} = \mathcal{R}(\mathbf{U}^{\top})$, and $\mathcal{V} = (\mathcal{V}^{\perp})^{\perp} = \mathcal{R}(\mathbf{U}^{\top})^{\perp} = \mathcal{N}(\mathbf{U})$ (by Exercise 3.11 and Theorem 3.4). **3.14**

a. Let $x \in \mathcal{V}$. Then, x = Px + (I - P)x. Note that $Px \in \mathcal{V}$, and $(I - P)x \in \mathcal{V}^{\perp}$. Therefore, x = Px + (I - P)x is an orthogonal decomposition of x with respect to \mathcal{V} . However, x = x + 0 is also an orthogonal decomposition of x with respect to \mathcal{V} . Since the orthogonal decomposition is unique, we must have x = Px.

b. Suppose P is an orthogonal projector onto \mathcal{V} . Clearly, $\mathcal{R}(P) \subset \mathcal{V}$ by definition. However, from part a, x = Px for all $x \in \mathcal{V}$, and hence $\mathcal{V} \subset \mathcal{R}(P)$. Therefore, $\mathcal{R}(P) = \mathcal{V}$.

3.15

To answer the question, we have to represent the quadratic form with a symmetric matrix as

$$oldsymbol{x}^{ op} \left(rac{1}{2} egin{bmatrix} 1 & -8 \ 1 & 1 \end{bmatrix} + rac{1}{2} egin{bmatrix} 1 & 1 \ -8 & 1 \end{bmatrix}
ight) oldsymbol{x} = oldsymbol{x}^{ op} egin{bmatrix} 1 & -7/2 \ -7/2 & 1 \end{bmatrix} oldsymbol{x}$$

The leading principal minors are $\Delta_1 = 1$ and $\Delta_2 = -45/4$. Therefore, the quadratic form is indefinite. **3.16**

The leading principal minors are $\Delta_1 = 2$, $\Delta_2 = 0$, $\Delta_3 = 0$, which are all nonnegative. However, the eigenvalues of \boldsymbol{A} are 0, -1.4641, 5.4641 (for example, use Matlab to quickly check this). This implies that the matrix \boldsymbol{A} is indefinite (by Theorem 3.7). An alternative way to show that \boldsymbol{A} is not positive semidefinite is to find a vector \boldsymbol{x} such that $\boldsymbol{x}^{\top} \boldsymbol{A} \boldsymbol{x} < 0$. So, let \boldsymbol{x} be an eigenvector of \boldsymbol{A} corresponding to its negative eigenvalue $\lambda = -1.4641$. Then, $\boldsymbol{x}^{\top} \boldsymbol{A} \boldsymbol{x} = \boldsymbol{x}^{\top} (\lambda \boldsymbol{x}) = \lambda \boldsymbol{x}^{\top} \boldsymbol{x} = \lambda \|\boldsymbol{x}\|^2 < 0$. For this example, we can take $\boldsymbol{x} = [0.3251, 0.3251, -0.8881]^{\top}$, for which we can verify that $\boldsymbol{x}^{\top} \boldsymbol{A} \boldsymbol{x} = -1.4643$.

3.17

a. The matrix \boldsymbol{Q} is indefinite, since $\Delta_2 = -1$ and $\Delta_3 = 2$.

b. Let $x \in M$. Then, $x_2 + x_3 = -x_1$, $x_1 + x_3 = -x_2$, and $x_1 + x_2 = -x_3$. Therefore,

$$\boldsymbol{x}^{\top}\boldsymbol{Q}\boldsymbol{x} = x_1(x_2 + x_3) + x_2(x_1 + x_3) + x_3(x_1 + x_2) = -(x_1^2 + x_2^2 + x_3^2).$$

This implies that the matrix Q is negative definite on the subspace \mathcal{M} .

3.18 _

a. We have

$$f(x_1, x_2, x_3) = x_2^2 = \begin{bmatrix} x_1, x_2, x_3 \end{bmatrix} \begin{bmatrix} 0 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 0 \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \\ x_3 \end{bmatrix}$$

Then,

$$\boldsymbol{Q} = \begin{bmatrix} 0 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 0 \end{bmatrix}$$

and the eigenvalues of Q are $\lambda_1 = 0$, $\lambda_2 = 1$, and $\lambda_3 = 0$. Therefore, the quadratic form is positive semidefinite.

b. We have

$$f(x_1, x_2, x_3) = x_1^2 + 2x_2^2 - x_1x_3 = \begin{bmatrix} x_1, x_2, x_3 \end{bmatrix} \begin{bmatrix} 1 & 0 & -\frac{1}{2} \\ 0 & 2 & 0 \\ -\frac{1}{2} & 0 & 0 \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \\ x_3 \end{bmatrix}$$

Then,

$$\boldsymbol{Q} = \begin{bmatrix} 1 & 0 & -\frac{1}{2} \\ 0 & 2 & 0 \\ -\frac{1}{2} & 0 & 0 \end{bmatrix}$$

and the eigenvalues of Q are $\lambda_1 = 2$, $\lambda_2 = (1 - \sqrt{2})/2$, and $\lambda_3 = (1 + \sqrt{2})/2$. Therefore, the quadratic form is indefinite.

c. We have

$$f(x_1, x_2, x_3) = x_1^2 + x_3^2 + 2x_1x_2 + 2x_1x_3 + 2x_2x_3 = \begin{bmatrix} x_1, x_2, x_3 \end{bmatrix} \begin{bmatrix} 1 & 1 & 1 \\ 1 & 0 & 1 \\ 1 & 1 & 1 \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \\ x_3 \end{bmatrix}.$$

Then,

$$oldsymbol{Q} = egin{bmatrix} 1 & 1 & 1 \ 1 & 0 & 1 \ 1 & 1 & 1 \end{bmatrix}$$

and the eigenvalues of Q are $\lambda_1 = 0$, $\lambda_2 = 1 - \sqrt{3}$, and $\lambda_3 = 1 + \sqrt{3}$. Therefore, the quadratic form is indefinite.

3.19 We have

$$f(x_1, x_2, x_3) = 4x_1^2 + x_2^2 + 9x_3^2 - 4x_1x_2 - 6x_2x_3 + 12x_1x_3$$

= $[x_1, x_2, x_3] \begin{bmatrix} 4 & -2 & 6 \\ -2 & 1 & -3 \\ 6 & -3 & 9 \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \\ x_3 \end{bmatrix}.$

Let

$$\boldsymbol{Q} = \begin{bmatrix} 4 & -2 & 6 \\ -2 & 1 & -3 \\ 6 & -3 & 9 \end{bmatrix}, \qquad \boldsymbol{x} = \begin{bmatrix} x_1 \\ x_2 \\ x_3 \end{bmatrix} = x_1 \boldsymbol{e}_1 + x_2 \boldsymbol{e}_2 + x_3 \boldsymbol{e}_3,$$

where e_1 , e_2 , and e_3 form the natural basis for \mathbb{R}^3 . Let v_1 , v_2 , and v_3 be another basis for \mathbb{R}^3 . Then, the vector \boldsymbol{x} is represented in the new basis as $\tilde{\boldsymbol{x}}$, where

 $\begin{aligned} \boldsymbol{x} &= [\boldsymbol{v}_1, \boldsymbol{v}_2, \boldsymbol{v}_3] \tilde{\boldsymbol{x}} = \boldsymbol{V} \tilde{\boldsymbol{x}}. \\ \text{Now, } f(\boldsymbol{x}) &= \boldsymbol{x}^\top \boldsymbol{Q} \boldsymbol{x} = (\boldsymbol{V} \tilde{\boldsymbol{x}})^\top \boldsymbol{Q} (\boldsymbol{V} \tilde{\boldsymbol{x}}) = \tilde{\boldsymbol{x}}^\top (\boldsymbol{V}^\top \boldsymbol{Q} \boldsymbol{V}) \tilde{\boldsymbol{x}} = \tilde{\boldsymbol{x}}^\top \tilde{\boldsymbol{Q}} \tilde{\boldsymbol{x}}, \text{ where} \end{aligned}$

$$\tilde{\boldsymbol{Q}} = \begin{bmatrix} \tilde{q}_{11} & \tilde{q}_{12} & \tilde{q}_{13} \\ \tilde{q}_{21} & \tilde{q}_{22} & \tilde{q}_{23} \\ \tilde{q}_{31} & \tilde{q}_{32} & \tilde{q}_{33} \end{bmatrix}$$

and $\tilde{q}_{ij} = \boldsymbol{v}_i \boldsymbol{Q} \boldsymbol{v}_j$ for i, j = 1, 2, 3. We will find a basis $\{\boldsymbol{v}_1, \boldsymbol{v}_2, \boldsymbol{v}_3\}$ such that $\tilde{q}_{ij} = 0$ for $i \neq j$, and is of the form

$$v_1 = \alpha_{11}e_1$$

$$v_2 = \alpha_{21}e_1 + \alpha_{22}e_2$$

$$v_3 = \alpha_{31}e_1 + \alpha_{32}e_2 + \alpha_{33}e_3$$

Because

$$\tilde{q}_{ij} = \boldsymbol{v}_i \boldsymbol{Q} \boldsymbol{v}_j = \boldsymbol{v}_i \boldsymbol{Q} (\alpha_{j1} \boldsymbol{e}_1 + \ldots + \alpha_{jj} \boldsymbol{e}_j) = \alpha_{j1} (\boldsymbol{v}_i \boldsymbol{Q} \boldsymbol{e}_1) + \ldots + \alpha_{jj} (\boldsymbol{v}_i \boldsymbol{Q} \boldsymbol{e}_j),$$

we deduce that if $v_i Q e_j = 0$ for j < i, then $v_i Q v_j = 0$. In this case,

$$\tilde{q}_{ii} = \boldsymbol{v}_i \boldsymbol{Q} \boldsymbol{v}_i = \boldsymbol{v}_i \boldsymbol{Q} (\alpha_{i1} \boldsymbol{e}_1 + \ldots + \alpha_{ii} \boldsymbol{e}_i) = \alpha_{i1} (\boldsymbol{v}_i \boldsymbol{Q} \boldsymbol{e}_1) + \ldots + \alpha_{ii} (\boldsymbol{v}_i \boldsymbol{Q} \boldsymbol{e}_i) = \alpha_{ii} (\boldsymbol{v}_i \boldsymbol{Q} \boldsymbol{e}_i)$$

Our task therefore is to find v_i (i = 1, 2, 3) such that

$$\begin{aligned} & \boldsymbol{v}_i \boldsymbol{Q} \boldsymbol{e}_j &= 0, \qquad j < \\ & \boldsymbol{v}_i \boldsymbol{Q} \boldsymbol{e}_i &= 1, \end{aligned}$$

i

and, in this case, we get

$$\tilde{\boldsymbol{Q}} = \begin{bmatrix} \alpha_{11} & 0 & 0 \\ 0 & \alpha_{22} & 0 \\ 0 & 0 & \alpha_{33} \end{bmatrix}.$$

• Case of i = 1.

From $\boldsymbol{v}_1^\top \boldsymbol{Q} \boldsymbol{e}_1 = 1$,

$$(\alpha_{11}\boldsymbol{e}_1)^{\top}\boldsymbol{Q}\boldsymbol{e}_1 = \alpha_{11}(\boldsymbol{e}_1^{\top}\boldsymbol{Q}\boldsymbol{e}_1) = \alpha_{11}q_{11} = 1.$$

Therefore,

$$\alpha_{11} = \frac{1}{q_{11}} = \frac{1}{\Delta_1} = \frac{1}{4} \qquad \Rightarrow \qquad \boldsymbol{v}_1 = \alpha_{11}\boldsymbol{e}_1 = \begin{bmatrix} \frac{1}{4} \\ 0 \\ 0 \\ 0 \end{bmatrix}.$$

• Case of i = 2.

From $\boldsymbol{v}_2^\top \boldsymbol{Q} \boldsymbol{e}_1 = 0$,

$$(\alpha_{21}\boldsymbol{e}_1 + \alpha_{22}\boldsymbol{e}_2)^{\top}\boldsymbol{Q}\boldsymbol{e}_1 = \alpha_{21}(\boldsymbol{e}_1^{\top}\boldsymbol{Q}\boldsymbol{e}_1) + \alpha_{22}(\boldsymbol{e}_2^{\top}\boldsymbol{Q}\boldsymbol{e}_1) = \alpha_{21}q_{11} + \alpha_{22}q_{21} = 0$$

From $\boldsymbol{v}_2^\top \boldsymbol{Q} \boldsymbol{e}_2 = 1$,

$$(\alpha_{21}\boldsymbol{e}_1 + \alpha_{22}\boldsymbol{e}_2)^{\top}\boldsymbol{Q}\boldsymbol{e}_2 = \alpha_{21}(\boldsymbol{e}_1^{\top}\boldsymbol{Q}\boldsymbol{e}_2) + \alpha_{22}(\boldsymbol{e}_2^{\top}\boldsymbol{Q}\boldsymbol{e}_2) = \alpha_{21}q_{12} + \alpha_{22}q_{22} = 1$$

Therefore,

$$\begin{bmatrix} q_{11} & q_{21} \\ q_{12} & q_{22} \end{bmatrix} \begin{bmatrix} \alpha_{21} \\ \alpha_{22} \end{bmatrix} = \begin{bmatrix} 0 \\ 1 \end{bmatrix}$$

But, since $\Delta_2 = 0$, this system of equations is inconsistent. Hence, in this problem $v_2^\top Q e_2 = 0$ should be satisfied instead of $v_2^\top Q e_2 = 1$ so that the system can have a solution. In this case, the diagonal matrix becomes

$$\tilde{\boldsymbol{Q}} = \begin{bmatrix} \alpha_{11} & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & \alpha_{33} \end{bmatrix},$$

and the system of equations become

$$\begin{bmatrix} q_{11} & q_{21} \\ q_{12} & q_{22} \end{bmatrix} \begin{bmatrix} \alpha_{21} \\ \alpha_{22} \end{bmatrix} = \begin{bmatrix} 0 \\ 0 \end{bmatrix} \Rightarrow \begin{bmatrix} \alpha_{21} \\ \alpha_{22} \end{bmatrix} = \begin{bmatrix} \frac{1}{2} \\ 1 \end{bmatrix} \alpha_{22},$$

where α_{22} is an arbitrary real number. Thus,

$$\boldsymbol{v}_2 = \alpha_{21}\boldsymbol{e}_1 + \alpha_{22}\boldsymbol{e}_2 = \begin{bmatrix} \frac{1}{2} \\ 1 \\ 0 \end{bmatrix} \boldsymbol{a},$$

where a is an arbitrary real number.

• Case of i = 3.

Since in this case $\Delta_3 = \det(\mathbf{Q}) = 0$, we will have to apply the same reasoning of the previous case and use the condition $\mathbf{v}_3^\top \mathbf{Q} \mathbf{e}_3 = 0$ instead of $\mathbf{v}_3^\top \mathbf{Q} \mathbf{e}_3 = 1$. In this way the diagonal matrix becomes

$$\tilde{\boldsymbol{Q}} = \begin{bmatrix} \alpha_{11} & 0 & 0 \\ 0 & 0 & 0 \\ 0 & 0 & 0 \end{bmatrix}$$

Thus, from $\boldsymbol{v}_3^\top \boldsymbol{Q} \boldsymbol{e}_1 = 0$, $\boldsymbol{v}_3^\top \boldsymbol{Q} \boldsymbol{e}_2 = 0$ and $\boldsymbol{v}_3^\top \boldsymbol{Q} \boldsymbol{e}_3 = 0$,

$$\begin{bmatrix} q_{11} & q_{21} & q_{31} \\ q_{12} & q_{22} & q_{32} \\ q_{13} & q_{23} & q_{33} \end{bmatrix} \begin{bmatrix} \alpha_{31} \\ \alpha_{32} \\ \alpha_{33} \end{bmatrix} = \mathbf{Q} \begin{bmatrix} \alpha_{31} \\ \alpha_{32} \\ \alpha_{33} \end{bmatrix} = \begin{bmatrix} \mathbf{Q} \\ \alpha_{31} \\ \alpha_{32} \\ \alpha_{33} \end{bmatrix} = \begin{bmatrix} 0 \\ 0 \\ 0 \end{bmatrix}$$

Therefore,

$$\begin{bmatrix} \alpha_{31} \\ \alpha_{32} \\ \alpha_{33} \end{bmatrix} = \begin{bmatrix} \alpha_{31} \\ 2\alpha_{31} + 3\alpha_{33} \\ \alpha_{33} \end{bmatrix},$$

where α_{31} and α_{33} are arbitrary real numbers. Thus,

$$\boldsymbol{v}_3 = \alpha_{31}\boldsymbol{e}_1 + \alpha_{32}\boldsymbol{e}_2 + \alpha_{33}\boldsymbol{e}_3 = \begin{bmatrix} b\\ 2b+3c\\c \end{bmatrix},$$

where b and c are arbitrary real numbers.

Finally,

$$\mathbf{V} = [x_1, x_2, x_3] = \begin{bmatrix} \frac{1}{4} & \frac{a}{2} & b\\ 0 & a & 2b + 3c\\ 0 & 0 & c \end{bmatrix},$$

where a, b, and c are arbitrary real numbers.

3.20 _

We represent this quadratic form as $f(\boldsymbol{x}) = \boldsymbol{x}^{\top} \boldsymbol{Q} \boldsymbol{x}$, where

$$\boldsymbol{Q} = \begin{bmatrix} 1 & \xi & -1 \\ \xi & 1 & 2 \\ -1 & 2 & 5 \end{bmatrix}$$

The leading principal minors of Q are $\Delta_1 = 1$, $\Delta_2 = 1 - \xi^2$, $\Delta_3 = -5\xi^2 - 4\xi$. For the quadratic form to be positive definite, all the leading principal minors of Q must be positive. This is the case if and only if $\xi \in (-4/5, 0).$

3.21

The matrix $\boldsymbol{Q} = \boldsymbol{Q}^{\top} > 0$ can be represented as $\boldsymbol{Q} = \boldsymbol{Q}^{1/2} \boldsymbol{Q}^{1/2}$, where $\boldsymbol{Q}^{1/2} = (\boldsymbol{Q}^{1/2})^{\top} > 0$. 1. Now, $\langle \boldsymbol{x}, \boldsymbol{x} \rangle_Q = (\boldsymbol{Q}^{1/2} \boldsymbol{x})^{\top} (\boldsymbol{Q}^{1/2} \boldsymbol{x}) = \| \boldsymbol{Q}^{1/2} \boldsymbol{x} \|^2 \ge 0$, and

$$\langle \boldsymbol{x}, \boldsymbol{x} \rangle_Q = 0 \quad \Leftrightarrow \quad \| \boldsymbol{Q}^{1/2} \boldsymbol{x} \|^2 = 0$$

 $\Leftrightarrow \quad \boldsymbol{Q}^{1/2} \boldsymbol{x} = \mathbf{0}$
 $\Leftrightarrow \quad \boldsymbol{x} = \mathbf{0}$

since $Q^{1/2}$ is nonsingular. 2. $\langle \boldsymbol{x}, \boldsymbol{y} \rangle_Q = \boldsymbol{x}^\top \boldsymbol{Q} \boldsymbol{y} = \boldsymbol{y}^\top \boldsymbol{Q}^\top \boldsymbol{x} = \boldsymbol{y}^\top \boldsymbol{Q} \boldsymbol{x} = \langle \boldsymbol{y}, \boldsymbol{x} \rangle_Q.$ 3. We have $\langle \boldsymbol{x}+\boldsymbol{y}, \boldsymbol{z}
angle_Q = (\boldsymbol{x}+\boldsymbol{y})^\top \boldsymbol{Q} \boldsymbol{z}$

$$egin{array}{rl} &=& oldsymbol{x}^{ op}oldsymbol{Q}oldsymbol{z}+oldsymbol{y}^{ op}oldsymbol{Q}oldsymbol{z}\ &=& \langleoldsymbol{x},oldsymbol{z}
angle_{Q}+\langleoldsymbol{y},oldsymbol{z}
angle_{Q}. \end{array}$$

4. $\langle r\boldsymbol{x}, \boldsymbol{y} \rangle_Q = (r\boldsymbol{x})^\top \boldsymbol{Q} \boldsymbol{y} = r\boldsymbol{x}^\top \boldsymbol{Q} \boldsymbol{y} = r \langle \boldsymbol{x}, \boldsymbol{y} \rangle_Q.$ 3.22 ____ We have

$$\|\boldsymbol{A}\|_{\infty} = \max\{\|\boldsymbol{A}\boldsymbol{x}\|_{\infty} : \|\boldsymbol{x}\|_{\infty} = 1\}.$$

We first show that $\|A\|_{\infty} \leq \max_{i} \sum_{k=1}^{n} |a_{ik}|$. For this, note that for each x such that $\|x\|_{\infty} = 1$, we have

$$\begin{aligned} \|\boldsymbol{A}\boldsymbol{x}\|_{\infty} &= \max_{i} \left| \sum_{k=1}^{n} a_{ik} x_{k} \right| \\ &\leq \max_{i} \sum_{k=1}^{n} |a_{ik}| |x_{k}| \\ &\leq \max_{i} \sum_{k=1}^{n} |a_{ik}|, \end{aligned}$$

since $|x_k| \leq \max_k |x_k| = ||\boldsymbol{x}||_{\infty} = 1$. Therefore,

$$\|\boldsymbol{A}\|_{\infty} \leq \max_{i} \sum_{k=1}^{n} |a_{ik}|.$$

To show that $\|\boldsymbol{A}\|_{\infty} = \max_{i} \sum_{k=1}^{n} |a_{ik}|$, it remains to find a $\tilde{\boldsymbol{x}} \in \mathbb{R}^{n}$, $\|\tilde{\boldsymbol{x}}\|_{\infty} = 1$, such that $\|\boldsymbol{A}\tilde{\boldsymbol{x}}\|_{\infty} = \max_{i} \sum_{k=1}^{n} |a_{ik}|$. So, let j be such that

$$\sum_{k=1}^{n} |a_{jk}| = \max_{i} \sum_{k=1}^{n} |a_{ik}|.$$

Define \tilde{x} by

$$\tilde{x}_k = \begin{cases} |a_{jk}|/a_{jk} & \text{if } a_{jk} \neq 0\\ 1 & \text{otherwise} \end{cases}.$$

Clearly $\|\tilde{\boldsymbol{x}}\|_{\infty} = 1$. Furthermore, for $i \neq j$,

$$\left|\sum_{k=1}^{n} a_{ik} \tilde{x}_{k}\right| \leq \sum_{k=1}^{n} |a_{ik}| \leq \max_{i} \sum_{k=1}^{n} |a_{ik}| = \sum_{k=1}^{n} |a_{jk}|$$

and

$$\left|\sum_{k=1}^{n} a_{jk} \tilde{x}_k\right| = \sum_{k=1}^{n} |a_{jk}|.$$

Therefore,

$$\|\boldsymbol{A}\tilde{\boldsymbol{x}}\|_{\infty} = \max_{i} \left| \sum_{k=1}^{n} a_{ik}\tilde{x}_{k} \right| = \sum_{k=1}^{n} |a_{jk}| = \max_{i} \sum_{k=1}^{n} |a_{ik}|.$$

3.23

We have

$$\|A\|_1 = \max\{\|Ax\|_1 : \|x\|_1 = 1\}.$$

We first show that $\|\mathbf{A}\|_1 \leq \max_k \sum_{i=1}^m |a_{ik}|$. For this, note that for each \mathbf{x} such that $\|\mathbf{x}\|_1 = 1$, we have

$$\begin{split} \|\boldsymbol{A}\boldsymbol{x}\|_{1} &= \sum_{i=1}^{m} \left| \sum_{k=1}^{n} a_{ik} x_{k} \right| \\ &\leq \sum_{i=1}^{m} \sum_{k=1}^{n} |a_{ik}| |x_{k}| \\ &\leq \sum_{k=1}^{n} |x_{k}| \sum_{i=1}^{m} |a_{ik}| \\ &\leq \left(\max_{k} \sum_{i=1}^{m} |a_{ik}| \right) \sum_{k=1}^{n} |x_{k}| \\ &\leq \max_{k} \sum_{i=1}^{m} |a_{ik}|, \end{split}$$

since $\sum_{k=1}^{n} |x_k| = \|\boldsymbol{x}\|_1 = 1$. Therefore,

$$\|\boldsymbol{A}\|_{1} \leq \max_{k} \sum_{i=1}^{m} |a_{ik}|.$$

To show that $\|\mathbf{A}\|_1 = \max_k \sum_{i=1}^m |a_{ik}|$, it remains to find a $\tilde{\mathbf{x}} \in \mathbb{R}^m$, $\|\tilde{\mathbf{x}}\|_1 = 1$, such that $\|\mathbf{A}\tilde{\mathbf{x}}\|_1 = \max_k \sum_{i=1}^m |a_{ik}|$. So, let j be such that

$$\sum_{i=1}^{m} |a_{ij}| = \max_{k} \sum_{i=1}^{m} |a_{ik}|.$$

Define \tilde{x} by

$$\tilde{x}_k = \begin{cases} 1 & \text{if } k = j \\ 0 & \text{otherwise} \end{cases}$$

Clearly $\|\tilde{\boldsymbol{x}}\|_1 = 1$. Furthermore,

$$\|\boldsymbol{A}\tilde{\boldsymbol{x}}\|_{1} = \sum_{i=1}^{m} \left| \sum_{k=1}^{n} a_{ik} \tilde{x}_{k} \right| = \sum_{i=1}^{m} |a_{ij}| = \max_{k} \sum_{i=1}^{m} |a_{ik}|.$$

4. Concepts from Geometry

4.1

 \Rightarrow : Let $S = \{ \boldsymbol{x} : \boldsymbol{A}\boldsymbol{x} = \boldsymbol{b} \}$ be a linear variety. Let $\boldsymbol{x}, \boldsymbol{y} \in S$ and $\alpha \in \mathbb{R}$. Then,

$$\boldsymbol{A}(\alpha \boldsymbol{x} + (1-\alpha)\boldsymbol{y}) = \alpha \boldsymbol{A}\boldsymbol{x} + (1-\alpha)\boldsymbol{A}\boldsymbol{y} = \alpha \boldsymbol{b} + (1-\alpha)\boldsymbol{b} = \boldsymbol{b}.$$

Therefore, $\alpha \boldsymbol{x} + (1 - \alpha) \boldsymbol{y} \in S$.

 $\Leftarrow: \text{ If } S \text{ is empty, we are done. So, suppose } \boldsymbol{x}_0 \in S. \text{ Consider the set } S_0 = S - \boldsymbol{x}_0 = \{\boldsymbol{x} - \boldsymbol{x}_0 : \boldsymbol{x} \in S\}.$ Clearly, for all $\boldsymbol{x}, \boldsymbol{y} \in S_0$ and $\alpha \in \mathbb{R}$, we have $\alpha \boldsymbol{x} + (1 - \alpha)\boldsymbol{y} \in S_0$. Note that $\boldsymbol{0} \in S_0$. We claim that S_0 is a subspace. To see this, let $\boldsymbol{x}, \boldsymbol{y} \in S_0$, and $\alpha \in \mathbb{R}$. Then, $\alpha \boldsymbol{x} = \alpha \boldsymbol{x} + (1 - \alpha)\boldsymbol{0} \in S_0$. Furthermore, $\frac{1}{2}\boldsymbol{x} + \frac{1}{2}\boldsymbol{y} \in S_0$, and therefore $\boldsymbol{x} + \boldsymbol{y} \in S_0$ by the previous argument. Hence, S_0 is a subspace. Therefore, by Exercise 3.13, there exists \boldsymbol{A} such that $S_0 = \mathcal{N}(\boldsymbol{A}) = \{\boldsymbol{x} : \boldsymbol{A}\boldsymbol{x} = \boldsymbol{0}\}.$ Define $\boldsymbol{b} = \boldsymbol{A}\boldsymbol{x}_0$. Then,

$$S = S_0 + x_0 = \{y + x_0 : y \in \mathcal{N}(A)\}$$

= $\{y + x_0 : Ay = 0\}$
= $\{y + x_0 : A(y + x_0) = b\}$
= $\{x : Ax = b\}.$

4.2 ____

Let $\boldsymbol{u}, \boldsymbol{v} \in \Theta = \{\boldsymbol{x} \in \mathbb{R}^n : \|\boldsymbol{x}\| \leq r\}$, and $\alpha \in [0, 1]$. Suppose $\boldsymbol{z} = \alpha \boldsymbol{u} + (1 - \alpha)\boldsymbol{v}$. To show that Θ is convex, we need to show that $\boldsymbol{z} \in \Theta$, i.e., $\|\boldsymbol{z}\| \leq r$. To this end,

$$\|\boldsymbol{z}\|^2 = (\alpha \boldsymbol{u}^\top + (1-\alpha)\boldsymbol{v}^\top)(\alpha \boldsymbol{u} + (1-\alpha)\boldsymbol{v})$$

= $\alpha^2 \|\boldsymbol{u}\|^2 + 2\alpha(1-\alpha)\boldsymbol{u}^\top \boldsymbol{v} + (1-\alpha)^2 \|\boldsymbol{v}\|^2.$

Since $\boldsymbol{u}, \boldsymbol{v} \in \Theta$, then $\|\boldsymbol{u}\|^2 \leq r^2$ and $\|\boldsymbol{v}\|^2 \leq r^2$. Furthermore, by the Cauchy-Schwarz Inequality, we have $\boldsymbol{u}^\top \boldsymbol{v} \leq \|\boldsymbol{u}\| \|\boldsymbol{v}\| \leq r^2$. Therefore,

$$\|\boldsymbol{z}\|^{2} \leq \alpha^{2} r^{2} + 2\alpha (1-\alpha)r^{2} + (1-\alpha)^{2}r^{2} = r^{2}.$$

Hence, $z \in \Theta$, which implies that Θ is a convex set, i.e., the any point on the line segment joining u and v is also in Θ .

 4.3_{-}

Let $\boldsymbol{u}, \boldsymbol{v} \in \Theta = \{\boldsymbol{x} \in \mathbb{R}^n : \boldsymbol{A}\boldsymbol{x} = \boldsymbol{b}\}$, and $\alpha \in [0, 1]$. Suppose $\boldsymbol{z} = \alpha \boldsymbol{u} + (1 - \alpha)\boldsymbol{v}$. To show that Θ is convex, we need to show that $\boldsymbol{z} \in \Theta$, i.e., $\boldsymbol{A}\boldsymbol{z} = \boldsymbol{b}$. To this end,

$$\begin{aligned} \boldsymbol{A}\boldsymbol{z} &= \boldsymbol{A}(\alpha\boldsymbol{u} + (1-\alpha)\boldsymbol{v}) \\ &= \alpha\boldsymbol{A}\boldsymbol{u} + (1-\alpha)\boldsymbol{A}\boldsymbol{v}. \end{aligned}$$

Since $u, v \in \Theta$, then Au = b and Av = b. Therefore,

$$Az = \alpha b + (1 - \alpha)b = b,$$

and hence $\boldsymbol{z} \in \Theta$.

4.4 _

Let $\boldsymbol{u}, \boldsymbol{v} \in \Theta = \{\boldsymbol{x} \in \mathbb{R}^n : \boldsymbol{x} \ge \boldsymbol{0}\}$, and $\alpha \in [0, 1]$. Suppose $\boldsymbol{z} = \alpha \boldsymbol{u} + (1 - \alpha)\boldsymbol{v}$. To show that Θ is convex, we need to show that $\boldsymbol{z} \in \Theta$, i.e., $\boldsymbol{z} \ge \boldsymbol{0}$. To this end, write $\boldsymbol{x} = [x_1, \dots, x_n]^\top$, $\boldsymbol{y} = [y_1, \dots, y_n]^\top$, and $\boldsymbol{z} = [z_1, \dots, z_n]^\top$. Then, $z_i = \alpha x_i + (1 - \alpha)y_i$, $i = 1, \dots, n$. Since $x_i, y_i \ge 0$, and $\alpha, 1 - \alpha \ge 0$, we have $z_i \ge 0$. Therefore, $\boldsymbol{z} \ge \boldsymbol{0}$, and hence $\boldsymbol{z} \in \Theta$.

5. Elements of Calculus

5.1 .

Observe that

$$\|\boldsymbol{A}^k\| \leq \|\boldsymbol{A}^{k-1}\|\|\boldsymbol{A}\| \leq \|\boldsymbol{A}^{k-2}\|\|\boldsymbol{A}\|^2 \leq \cdots \leq \|\boldsymbol{A}\|^k.$$

Therefore, if $\|\mathbf{A}\| < 1$, then $\lim_{k \to \infty} \|\mathbf{A}^k\| = \mathbf{O}$ which implies that $\lim_{k \to \infty} \mathbf{A}^k = \mathbf{O}$.

5.2 _

For the case when A has all real eigenvalues, the proof is simple. Let λ be the eigenvalue of A with largest absolute value, and x the corresponding (normalized) eigenvector, i.e., $Ax = \lambda x$ and ||x|| = 1. Then,

$$\|\boldsymbol{A}\| \geq \|\boldsymbol{A}\boldsymbol{x}\| = \|\lambda\boldsymbol{x}\| = |\lambda|\|\boldsymbol{x}\| = |\lambda|,$$

which completes the proof for this case.

In general, the eigenvalues of A and the corresponding eigenvectors may be complex. In this case, we proceed as follows (see [41]). Consider the matrix

$$B = rac{A}{\|A\| + arepsilon},$$

where ε is a positive real number. We have

$$\|\boldsymbol{B}\| = \frac{\|\boldsymbol{A}\|}{\|\boldsymbol{A}\| + \varepsilon} < 1.$$

By Exercise 5.1, $\mathbf{B}^k \to \mathbf{O}$ as $k \to \infty$, and thus by Lemma 5.1, $|\lambda_i(\mathbf{B})| < 1, i = 1, ..., n$. On the other hand, for each i = 1, ..., n,

$$\lambda_i(\boldsymbol{B}) = rac{\lambda_i(\boldsymbol{A})}{\|\boldsymbol{A}\| + \varepsilon},$$

and thus

$$|\lambda_i(\boldsymbol{B})| = rac{|\lambda_i(\boldsymbol{A})|}{\|\boldsymbol{A}\| + \varepsilon} < 1.$$

 $|\lambda_i(\boldsymbol{A})| < \|\boldsymbol{A}\| + \varepsilon.$

which gives

Since the above arguments hold for any $\varepsilon > 0$, we have $|\lambda_i(\mathbf{A})| \le ||\mathbf{A}||$. 5.3

a.
$$\nabla f(\boldsymbol{x}) = (\boldsymbol{a}\boldsymbol{b}^\top + \boldsymbol{b}\boldsymbol{a}^\top)\boldsymbol{x}.$$

b. $F(\boldsymbol{x}) = \boldsymbol{a}\boldsymbol{b}^\top + \boldsymbol{b}\boldsymbol{a}^\top.$

5.4 ______ We have

 $\quad \text{and} \quad$

$$Df(\mathbf{x}) = [x_1/3, x_2/2],$$

$$\frac{d\boldsymbol{g}}{dt}(t) = \begin{bmatrix} 3\\2 \end{bmatrix}.$$

By the chain rule,

$$\begin{aligned} \frac{d}{dt}F(t) &= Df(\boldsymbol{g}(t))\frac{d\boldsymbol{g}}{dt}(t) \\ &= \left[(3t+5)/3,(2t-6)/2\right] \begin{bmatrix} 3\\ 2 \end{bmatrix} \\ &= 5t-1. \end{aligned}$$

5.5 ____

We have

 $\quad \text{and} \quad$

$$\frac{\partial \boldsymbol{g}}{\partial s}(s,t) = \begin{bmatrix} 4\\2 \end{bmatrix}, \qquad \qquad \frac{\partial \boldsymbol{g}}{\partial t}(s,t) = \begin{bmatrix} 3\\1 \end{bmatrix}.$$

 $Df(\mathbf{x}) = [x_2/2, x_1/2],$

By the chain rule,

$$\begin{aligned} \frac{\partial}{\partial s} f(\boldsymbol{g}(s,t)) &= Df(\boldsymbol{g}(t)) \frac{\partial \boldsymbol{g}}{\partial s}(s,t) \\ &= \frac{1}{2} [2s+t, 4s+3t] \begin{bmatrix} 4\\ 2 \end{bmatrix} \\ &= 8s+5t, \end{aligned}$$

 $\quad \text{and} \quad$

$$\frac{\partial}{\partial t} f(\boldsymbol{g}(s,t)) = Df(\boldsymbol{g}(t)) \frac{\partial \boldsymbol{g}}{\partial t}(s,t)$$

$$= \frac{1}{2} [2s+t, 4s+3t] \begin{bmatrix} 3\\1 \end{bmatrix}$$

$$= 5s+3t.$$

5.6 ____

We have

and

$$Df(\boldsymbol{x}) = [3x_1^2x_2x_3^2 + x_2, \ x_1^3x_3^2 + x_1, \ 2x_1^3x_2x_3 + 1]$$

$$\frac{d\boldsymbol{x}}{dt}(t) = \begin{bmatrix} e^t + 3t^2 \\ 2t \\ 1 \end{bmatrix}.$$