

- [3] Bliss, C. I., "Statistics in Biology," *Statistical Methods for Research in the Natural Sciences*, vol. 2, New York: McGraw-Hill, 1967.
- [4] Cleveland, W. S., and D. A. Relles, "Clustering by Identification with Special Application to Two-Way Tables of Counts," *Journal of the American Statistical Association*, 70, no. 351 (1975), 626-630.
- [5] Cochran, W. G., *Sampling Techniques* (3rd ed.), New York: John Wiley, 1977.
- [6] Cochran, W. G., and G. M. Cox, *Experimental Designs* (2nd ed.), New York: John Wiley, 1957.
- [7] Davis, J. C., "Information Contained in Sediment Size Analysis," *Mathematical Geology*, 2, no. 2 (1970), 105-112.
- [8] Dunham, R. B., and D. J. Kravetz, "Canonical Correlation Analysis in a Predictive System," *Journal of Experimental Education*, 43, no. 4 (1975), 35-42.
- [9] Fritts, H. C., T. J. Blasing, B. P. Hayden, and J. E. Kutzbach, "Multivariate Techniques for Specifying Tree-Growth and Climate Relationships and Reconstructing Anomalies in Paleoclimate," *Journal of Applied Meteorology*, 10, no. 5 (1971), 845-864.
- [10] Halinar, J. C., "Principal Component Analysis in Plant Breeding," unpublished report based on data collected by Dr. F. A. Bliss, University of Wisconsin, 1970.
- [11] Klatzky, S. R., and R. W. Hodge, "A Canonical Correlation Analysis of Occupational Mobility," *Journal of the American Statistical Association*, 66, no. 333 (1971), 16-22.
- [12] Linden, M., "Factor Analytic Study of Olympic Decathlon Data," *Research Quarterly*, 48, no. 3 (Oct. 1977), 562-568.
- [13] Mather, P. M., "Study of Factors Influencing Variation in Size Characteristics in Fluvio-glacial Sediments," *Mathematical Geology*, 4, no. 3 (1972), 219-234.
- [14] Marriott, F. H. C., *The Interpretation of Multiple Observations*, London: Academic Press, 1974.
- [15] Spinner, K. I., "From Generation to Generation: The Transmission of Occupation," Ph.D. dissertation, University of Wisconsin, 1977.
- [16] Timm, N. H., *Multivariate Analysis with Applications in Education and Psychology*, Monterey, Calif.: Brooks/Cole, 1975.
- [17] Trieschmann, J. S., and G. E. Pinches, "A Multivariate Model for Predicting Financially Distressed P-L Insurers," *Journal of Risk and Insurance*, 40, no. 3 (1973), 327-338.
- [18] Tukey, J. W., *Exploratory Data Analysis*, Reading, Mass.: Addison-Wesley, 1977.

# 2

## Matrix Algebra and Random Vectors

### 2.1 INTRODUCTION

We saw in Chapter 1 that multivariate data can be conveniently displayed as an array of numbers. In general, a rectangular array of numbers with, for instance,  $p$  rows and  $n$  columns is called a *matrix* of dimension  $p \times n$ . The study of multivariate methods is greatly facilitated by the use of matrix algebra.

The matrix algebra results presented in this chapter will enable us to concisely state statistical models. Moreover, the formal relations expressed in matrix terms are easily programmed on computers to allow the routine calculation of important statistical quantities.

We begin by introducing some very basic concepts that are essential to both our geometrical interpretations and algebraic explanations of subsequent statistical techniques. If you have not been previously exposed to the rudiments of matrix algebra, you may prefer to follow the brief refresher in the next section by the more detailed review provided in Supplement 2A.

### 2.2 SOME BASICS OF MATRIX AND VECTOR ALGEBRA

#### Vectors

An array  $x$  of  $n$  real numbers  $x_1, x_2, \dots, x_n$  is called a *vector* and it is written as



$$\mathbf{x} = \begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_n \end{bmatrix} \quad \text{or} \quad \mathbf{x}' = [x_1, x_2, \dots, x_n]$$

where the prime denotes the operation of *transposing* a column to a row. A vector  $\mathbf{x}$  can be represented geometrically as a directed line in  $n$ -dimensions with component  $x_1$  along the first axis,  $x_2$  along the second axis,  $\dots$ ,  $x_n$  along the  $n$ th axis. This is illustrated in Figure 2.1 for  $n = 3$ .

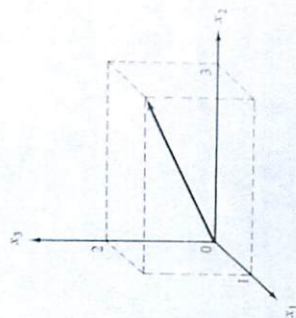


Figure 2.1 The vector  $\mathbf{x}' = [1, 2, 3]$

A vector can be *expanded* or *contracted* by multiplying by a constant  $c$ . In particular, we define the vector  $c\mathbf{x}$  as

$$c\mathbf{x} = \begin{bmatrix} cx_1 \\ cx_2 \\ \vdots \\ cx_n \end{bmatrix}$$

That is,  $c\mathbf{x}$  is the vector obtained by multiplying each element of  $\mathbf{x}$  by  $c$ . [See Figure 2.2(a).]

Two vectors may be added. Addition of  $\mathbf{x}$  and  $\mathbf{y}$  is defined as

$$\mathbf{x} + \mathbf{y} = \begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_n \end{bmatrix} + \begin{bmatrix} y_1 \\ y_2 \\ \vdots \\ y_n \end{bmatrix} = \begin{bmatrix} x_1 + y_1 \\ x_2 + y_2 \\ \vdots \\ x_n + y_n \end{bmatrix}$$

so that  $\mathbf{x} + \mathbf{y}$  is the vector with  $i$ th element  $x_i + y_i$ .

The sum of two vectors emanating from the origin is the diagonal of the parallelogram formed with the two original vectors as adjacent sides. This geometrical interpretation is illustrated in Figure 2.2(b).

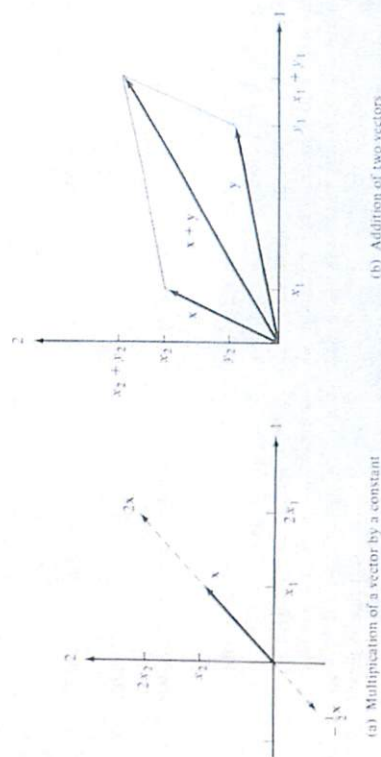


Figure 2.2 (a) Multiplication of a vector by a constant  
(b) Addition of two vectors

A vector has both direction and length. In  $n = 2$  dimensions, we consider the vector

$$\mathbf{x} = \begin{bmatrix} x_1 \\ x_2 \end{bmatrix}$$

The length of  $\mathbf{x}$ , written  $L_x$ , is defined to be

$$L_x = \sqrt{x_1^2 + x_2^2}$$

Geometrically, the length of a vector in two dimensions can be viewed as the hypotenuse of a right triangle. This is demonstrated schematically in Figure 2.3. The length of a vector  $\mathbf{x}' = [x_1, x_2, \dots, x_n]$ , with  $n$  components, is defined by

$$L_x = \sqrt{x_1^2 + x_2^2 + \dots + x_n^2} \quad (2-1)$$

Multiplication of a vector  $\mathbf{x}$  by a scalar  $c$  changes the length. From Equation (2-1),

$$\begin{aligned} L_{c\mathbf{x}} &= \sqrt{c^2 x_1^2 + c^2 x_2^2 + \dots + c^2 x_n^2} \\ &= |c| \sqrt{x_1^2 + x_2^2 + \dots + x_n^2} = |c| L_x \end{aligned}$$

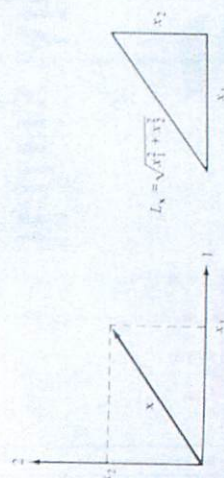


Figure 2.3 Length of  $\mathbf{x} = \sqrt{x_1^2 + x_2^2}$

Multiplication by  $c$  does not change the direction of the vector  $\mathbf{x}$  if  $c > 0$ . However, a negative value of  $c$  creates a vector with a direction opposite that of  $\mathbf{x}$ . From

$$L_{c\mathbf{x}} = |c| L_{\mathbf{x}} \quad (2-2)$$

it is clear that  $\mathbf{x}$  is expanded if  $|c| > 1$  and contracted if  $0 < |c| < 1$ . [Recall Figure 2.2(a).] Choosing  $c = L_{\mathbf{x}}^{-1}$  we obtain the unit vector  $L_{\mathbf{x}}^{-1}\mathbf{x}$ , which has length 1 and lies in the direction of  $\mathbf{x}$ .

A second geometrical concept is *angle*. Consider two vectors in a plane and the angle,  $\theta$ , between them, as in Figure 2.4. From Figure 2.4,  $\theta$  can be represented as the difference between the angles  $\theta_1$  and  $\theta_2$  formed by the two vectors and the first coordinate axis. Since, by definition,

$$\cos(\theta_1) = \frac{x_1}{L_{\mathbf{x}}} \quad \cos(\theta_2) = \frac{y_1}{L_{\mathbf{y}}}$$

$$\sin(\theta_1) = \frac{x_2}{L_{\mathbf{x}}} \quad \sin(\theta_2) = \frac{y_2}{L_{\mathbf{y}}}$$

and

$$\cos(\theta) = \cos(\theta_2 - \theta_1) = \cos(\theta_2)\cos(\theta_1) + \sin(\theta_2)\sin(\theta_1)$$

the angle  $\theta$  between the two vectors  $\mathbf{x}' = [x_1, x_2]$  and  $\mathbf{y}' = [y_1, y_2]$  is specified by

$$\cos(\theta) = \cos(\theta_2 - \theta_1) = \left(\frac{y_1}{L_{\mathbf{y}}}\right)\left(\frac{x_1}{L_{\mathbf{x}}}\right) + \left(\frac{y_2}{L_{\mathbf{y}}}\right)\left(\frac{x_2}{L_{\mathbf{x}}}\right) = \frac{x_1 y_1 + x_2 y_2}{L_{\mathbf{x}} L_{\mathbf{y}}} \quad (2-3)$$

We find it convenient to introduce the *inner product* of two vectors. For  $n = 2$  dimensions, the inner product of  $\mathbf{x}$  and  $\mathbf{y}$  is

$$\mathbf{x}'\mathbf{y} = x_1 y_1 + x_2 y_2$$

With this definition and Equation (2-3),

$$L_{\mathbf{x}} = \sqrt{\mathbf{x}'\mathbf{x}}, \quad \cos(\theta) = \frac{\mathbf{x}'\mathbf{y}}{L_{\mathbf{x}} L_{\mathbf{y}}} = \frac{\mathbf{x}'\mathbf{y}}{\sqrt{\mathbf{x}'\mathbf{x}} \sqrt{\mathbf{y}'\mathbf{y}}}$$

Since  $\cos(90^\circ) = \cos(270^\circ) = 0$  and  $\cos(\theta) = 0$  only if  $\mathbf{x}'\mathbf{y} = 0$ ,  $\mathbf{x}$  and  $\mathbf{y}$  are perpendicular when  $\mathbf{x}'\mathbf{y} = 0$ .

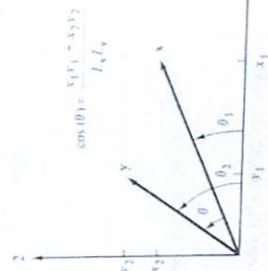


Figure 2.4 The angle  $\theta$  between  $\mathbf{x}$  and  $\mathbf{y}$  is  $\theta = \theta_2 - \theta_1$ .

For an arbitrary number of dimensions,  $n$ , we define the inner product of  $\mathbf{x}$  and  $\mathbf{y}$  as

$$\mathbf{x}'\mathbf{y} = x_1 y_1 + x_2 y_2 + \dots + x_n y_n \quad (2-4)$$

The inner product is denoted by either  $\mathbf{x}'\mathbf{y}$  or  $\mathbf{y}'\mathbf{x}$ .

Using the inner product, we have the natural extension of length and angle to vectors of  $n$  components:

$$L_{\mathbf{x}} = \text{length of } \mathbf{x} = \sqrt{\mathbf{x}'\mathbf{x}} \quad (2-5)$$

$$\cos(\theta) = \frac{\mathbf{x}'\mathbf{y}}{L_{\mathbf{x}} L_{\mathbf{y}}} = \frac{\mathbf{x}'\mathbf{y}}{\sqrt{\mathbf{x}'\mathbf{x}} \sqrt{\mathbf{y}'\mathbf{y}}} \quad (2-6)$$

Since, again,  $\cos(\theta) = 0$  only if  $\mathbf{x}'\mathbf{y} = 0$ , we say  $\mathbf{x}$  and  $\mathbf{y}$  are *perpendicular* when  $\mathbf{x}'\mathbf{y} = 0$ .

### Example 2.1

Given the vectors  $\mathbf{x}' = [1, 3, 2]$  and  $\mathbf{y}' = [-2, 1, -1]$ , find  $3\mathbf{x}$  and  $\mathbf{x} + \mathbf{y}$ . Next, determine the length of  $\mathbf{x}$ , the length of  $\mathbf{y}$ , and the angle between  $\mathbf{x}$  and  $\mathbf{y}$ . Also, check that the length of  $3\mathbf{x}$  is three times the length of  $\mathbf{x}$ .

First,

$$3\mathbf{x} = 3 \begin{bmatrix} 1 \\ 3 \\ 2 \end{bmatrix} = \begin{bmatrix} 3 \\ 9 \\ 6 \end{bmatrix}$$

$$\mathbf{x} + \mathbf{y} = \begin{bmatrix} 1 \\ 3 \\ 2 \end{bmatrix} + \begin{bmatrix} -2 \\ 1 \\ -1 \end{bmatrix} = \begin{bmatrix} -1 \\ 4 \\ 1 \end{bmatrix}$$

Next,  $\mathbf{x}'\mathbf{x} = 1^2 + 3^2 + 2^2 = 14$ ,  $\mathbf{y}'\mathbf{y} = (-2)^2 + 1^2 + (-1)^2 = 6$ , and  $\mathbf{x}'\mathbf{y} = 1(-2) + 3(1) + 2(-1) = -1$ . Therefore,

$$L_{\mathbf{x}} = \sqrt{\mathbf{x}'\mathbf{x}} = \sqrt{14} = 3.742 \quad L_{\mathbf{y}} = \sqrt{\mathbf{y}'\mathbf{y}} = \sqrt{6} = 2.449$$

and

$$\cos(\theta) = \frac{\mathbf{x}'\mathbf{y}}{L_{\mathbf{x}} L_{\mathbf{y}}} = \frac{-1}{3.742 \times 2.449} = -.109$$

so  $\theta = 96.3$  degrees. Finally,

$$L_{3\mathbf{x}} = \sqrt{3^2 + 9^2 + 6^2} = \sqrt{126} \quad \text{and} \quad 3L_{\mathbf{x}} = 3\sqrt{14} = \sqrt{126}$$

showing  $L_{3\mathbf{x}} = 3L_{\mathbf{x}}$ . ■

A pair of vectors  $\mathbf{x}$  and  $\mathbf{y}$  of the same dimension is said to be *linearly dependent* if there exist constants  $c_1$  and  $c_2$ , both not zero, such that

$$c_1 \mathbf{x} + c_2 \mathbf{y} = \mathbf{0}$$

A set of vectors  $\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_k$  is said to be *linearly dependent* if there exist constants  $c_1, c_2, \dots, c_k$ , not all zero, such that

$$c_1 \mathbf{x}_1 + c_2 \mathbf{x}_2 + \dots + c_k \mathbf{x}_k = \mathbf{0} \quad (2-7)$$



Linear dependence implies at least one vector in the set can be written as a linear combination of the other vectors. Vectors of the same dimension that are not linearly dependent are said to be *linearly independent*.

### Example 2.2

Consider the set of vectors

$$\mathbf{x}_1 = \begin{bmatrix} 1 \\ 2 \\ 1 \end{bmatrix}, \quad \mathbf{x}_2 = \begin{bmatrix} 1 \\ 0 \\ -1 \end{bmatrix}, \quad \mathbf{x}_3 = \begin{bmatrix} 1 \\ -2 \\ 1 \end{bmatrix}$$

Setting

$$c_1\mathbf{x}_1 + c_2\mathbf{x}_2 + c_3\mathbf{x}_3 = \mathbf{0}$$

implies

$$c_1 + c_2 + c_3 = 0$$

$$2c_1 - 2c_3 = 0$$

$$c_1 - c_2 + c_3 = 0$$

with the unique solution  $c_1 = c_2 = c_3 = 0$ . As we cannot find three constants  $c_1$ ,  $c_2$ , and  $c_3$ , not all zero, such that  $c_1\mathbf{x}_1 + c_2\mathbf{x}_2 + c_3\mathbf{x}_3 = \mathbf{0}$ , the vectors  $\mathbf{x}_1$ ,  $\mathbf{x}_2$ , and  $\mathbf{x}_3$  are *linearly independent*. ■

The *projection* (or shadow) of a vector  $\mathbf{x}$  on a vector  $\mathbf{y}$  is

$$\text{Projection of } \mathbf{x} \text{ on } \mathbf{y} = \frac{(\mathbf{x} \cdot \mathbf{y})}{\mathbf{y} \cdot \mathbf{y}} \mathbf{y} = \frac{(\mathbf{x} \cdot \mathbf{y})}{L_y} \frac{\mathbf{y}}{L_y} \quad (2-8)$$

where the vector  $L_y^{-1}\mathbf{y}$  has unit length. The *length* of the projection is

$$\text{Length of projection} = \left| \frac{\mathbf{x} \cdot \mathbf{y}}{L_y} \right| = L_x \left| \frac{\mathbf{x} \cdot \mathbf{y}}{L_x L_y} \right| = L_x |\cos(\theta)| \quad (2-9)$$

where  $\theta$  is the angle between  $\mathbf{x}$  and  $\mathbf{y}$  (see Figure 2.5).

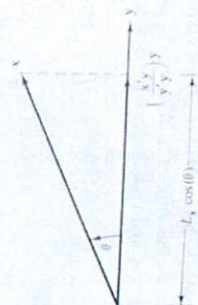


Figure 2.5 The projection of  $\mathbf{x}$  on  $\mathbf{y}$ .

### Matrices

A *matrix* is any rectangular array of real numbers. We denote an arbitrary array of  $p$  rows and  $n$  columns by

$$\mathbf{A} = \begin{bmatrix} a_{11} & a_{12} & \cdots & a_{1n} \\ a_{21} & a_{22} & \cdots & a_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{p1} & a_{p2} & \cdots & a_{pn} \end{bmatrix} \quad (p \times n)$$

Many of the vector concepts introduced above have direct generalizations to matrices.

The *transpose* operation,  $\mathbf{A}'$ , of a matrix changes the columns into rows so that the first column of  $\mathbf{A}$  becomes the first row of  $\mathbf{A}'$ , the second column becomes the second row, and so forth.

### Example 2.3

If

$$\mathbf{A} = \begin{bmatrix} 3 & -1 & 2 \\ 1 & 5 & 4 \end{bmatrix} \quad (2 \times 3)$$

then

$$\mathbf{A}' = \begin{bmatrix} 3 & 1 \\ -1 & 5 \\ 2 & 4 \end{bmatrix} \quad (3 \times 2)$$

A matrix may also be multiplied by a constant  $c$ . The product  $c\mathbf{A}$  is the matrix that results from multiplying each element of  $\mathbf{A}$  by  $c$ . Thus

$$c\mathbf{A} = \begin{bmatrix} ca_{11} & ca_{12} & \cdots & ca_{1n} \\ ca_{21} & ca_{22} & \cdots & ca_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ ca_{p1} & ca_{p2} & \cdots & ca_{pn} \end{bmatrix} \quad (p \times n)$$

Two matrices,  $\mathbf{A}$  and  $\mathbf{B}$ , of the same dimensions can be added. The sum  $\mathbf{A} + \mathbf{B}$  has  $(i, j)$ th entry  $a_{ij} + b_{ij}$ .

### Example 2.4

$$\text{If } \mathbf{A} = \begin{bmatrix} 0 & 3 & 1 \\ 1 & -1 & 1 \end{bmatrix} \quad (2 \times 3) \quad \text{and} \quad \mathbf{B} = \begin{bmatrix} 1 & -2 & -3 \\ 2 & 5 & 1 \end{bmatrix} \quad (2 \times 3)$$

$$\text{and} \quad 4\mathbf{A} = \begin{bmatrix} 0 & 12 & 4 \\ 4 & -4 & 4 \end{bmatrix} \quad (2 \times 3)$$

$$\mathbf{A} + \mathbf{B} = \begin{bmatrix} 0+1 & 3-2 & 1-3 \\ 1+2 & -1+5 & 1+1 \end{bmatrix} = \begin{bmatrix} 1 & 1 & -2 \\ 3 & 4 & 2 \end{bmatrix} \quad (2 \times 3)$$

It is also possible to define *matrix* multiplication if the dimensions of the matrices conform in the following manner. When  $\mathbf{A}$  is  $(p \times k)$  and  $\mathbf{B}$  is  $(k \times n)$ , so that the number of elements in a row of  $\mathbf{A}$  is the same as the number of

elements in a column of  $B$ , we can form the matrix product  $AB$ . An element of the new matrix  $AB$  is formed by taking the inner product of each row of  $A$  with each column of  $B$ .

The matrix product  $AB$  is

$A$   $B$  = the  $(p \times n)$  matrix whose entry in the  $i$ th row and  $j$ th column is the inner product of the  $i$ th row of  $A$  and  $j$ th column of  $B$

or

$$(i, j) \text{ entry of } AB = a_{i1}b_{1j} + a_{i2}b_{2j} + \cdots + a_{ik}b_{kj} = \sum_{l=1}^k a_{il}b_{lj} \quad (2-10)$$

When  $k = 4$ , we have four products to add for each entry in the matrix  $AB$ . Thus

$$A \quad B = \begin{bmatrix} a_{11} & a_{12} & a_{13} & a_{14} \\ \vdots & \vdots & \vdots & \vdots \\ a_{i1} & a_{i2} & a_{i3} & a_{i4} \\ \vdots & \vdots & \vdots & \vdots \\ a_{p1} & a_{p2} & a_{p3} & a_{p4} \end{bmatrix} \begin{bmatrix} b_{11} & \cdots & b_{1j} & \cdots & b_{1n} \\ b_{21} & \cdots & b_{2j} & \cdots & b_{2n} \\ b_{31} & \cdots & b_{3j} & \cdots & b_{3n} \\ b_{41} & \cdots & b_{4j} & \cdots & b_{4n} \end{bmatrix}$$

Column  
 $j$

$$= \text{Row } i \begin{bmatrix} \cdots (a_{i1}b_{1j} + a_{i2}b_{2j} + a_{i3}b_{3j} + a_{i4}b_{4j}) \cdots \end{bmatrix}$$

#### Example 2.5

If

$$A = \begin{bmatrix} 3 & -1 & 2 \\ 1 & 5 & 4 \end{bmatrix}, \quad B = \begin{bmatrix} -2 \\ 7 \\ 9 \end{bmatrix}, \quad \text{and} \quad C = \begin{bmatrix} 2 & 0 \\ 1 & -1 \end{bmatrix}$$

then

$$A \quad B = \begin{bmatrix} 3 & -1 & 2 \\ 1 & 5 & 4 \end{bmatrix} \begin{bmatrix} -2 \\ 7 \\ 9 \end{bmatrix} = \begin{bmatrix} 3(-2) + (-1)(7) + 2(9) \\ 1(-2) + 5(7) + 4(9) \end{bmatrix} = \begin{bmatrix} 5 \\ 69 \end{bmatrix}$$

$$C \quad A = \begin{bmatrix} 2 & 0 \\ 1 & -1 \end{bmatrix} \begin{bmatrix} 3 & -1 & 2 \\ 1 & 5 & 4 \end{bmatrix}$$

$$= \begin{bmatrix} 2(3) + 0(1) & 2(-1) + 0(5) & 2(2) + 0(4) \\ 1(3) - 1(1) & 1(-1) - 1(5) & 1(2) - 1(4) \end{bmatrix} = \begin{bmatrix} 6 & -2 & 4 \\ 2 & -6 & -2 \end{bmatrix}$$

(2 × 3)

Square matrices will be of special importance in our development of statistical methods. A square matrix is said to be *symmetric* if  $A = A'$  or  $a_{ij} = a_{ji}$  for all  $i$  and  $j$ .

#### Example 2.6

The matrix

$$\begin{bmatrix} 3 & 5 \\ 5 & -2 \end{bmatrix}$$

is symmetric; the matrix

$$\begin{bmatrix} 3 & 6 \\ 4 & -2 \end{bmatrix}$$

is not symmetric.

When two square matrices  $A$  and  $B$  are of the same dimension, both products  $AB$  and  $BA$  are defined, although they need not be equal (see Supplement 2A). If we let  $I$  denote the square matrix with ones on the diagonal and zeros elsewhere, it follows from the definition of matrix multiplication that the  $(i, j)$  entry of  $AI$  is  $a_{i1} \times 0 + \cdots + a_{i,j-1} \times 0 + a_{ij} \times 1 + a_{i,j+1} \times 0 + \cdots + a_{ik} \times 0 = a_{ij}$ , so  $AI = A$ . Similarly,  $IA = A$  and

$$\begin{matrix} I & A & A & I \\ (k \times k) & (k \times k) & (k \times k) & (k \times k) \end{matrix} = \begin{matrix} A & A & A \\ (k \times k) & (k \times k) & (k \times k) \end{matrix} \quad \text{for any } A \quad (2-11)$$

The matrix  $I$  acts like 1 in ordinary multiplication ( $1 \cdot a = a \cdot 1 = a$ ), so it is called the *identity* matrix.

The fundamental scalar relation about the existence of an inverse number  $a^{-1}$  such that  $a^{-1}a = aa^{-1} = 1$ , if  $a \neq 0$ , has the following matrix algebra extension. If there exists a matrix  $B$  such that

$$\begin{matrix} B & A & A & B \\ (k \times k) & (k \times k) & (k \times k) & (k \times k) \end{matrix} = \begin{matrix} I & I \\ (k \times k) & (k \times k) \end{matrix}$$

then  $B$  is called the *inverse* of  $A$  and is denoted by  $A^{-1}$ .

The technical condition that a  $k \times k$  inverse exists is that the  $k$  columns  $a_1, a_2, \dots, a_k$  of  $A$  are linearly independent. That is, the existence of  $A^{-1}$  is equivalent to

$$c_1a_1 + c_2a_2 + \cdots + c_ka_k = 0 \quad \text{only if} \quad c_1 = \cdots = c_k = 0 \quad (2-12)$$

(See Result 2A.9 in Supplement 2A.)

#### Example 2.7

For

$$A = \begin{bmatrix} 3 & 2 \\ 4 & 1 \end{bmatrix}$$



you may verify that

$$\begin{bmatrix} -.2 & .4 \\ .8 & -.6 \end{bmatrix} \begin{bmatrix} 3 & 2 \\ 4 & 1 \end{bmatrix} = \begin{bmatrix} (-.2)3 + (.4)4 & (-.2)2 + (.4)1 \\ (.8)3 + (-.6)4 & (.8)2 + (-.6)1 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}$$

so

$$\begin{bmatrix} -.2 & .4 \\ .8 & -.6 \end{bmatrix}$$

is  $A^{-1}$ . We note that

$$c_1 \begin{bmatrix} 3 \\ 4 \end{bmatrix} + c_2 \begin{bmatrix} 2 \\ 1 \end{bmatrix} = \begin{bmatrix} 0 \\ 0 \end{bmatrix}$$

implies  $c_1 = c_2 = 0$ , so the columns of  $A$  are linearly independent. This confirms the condition stated in (2-12). ■

A method for computing an inverse, when it exists, is given in Supplement 2A. The routine, but lengthy, calculations are usually relegated to a computer, especially when the dimension is greater than three. Even so, you must be forewarned that if the column sum in (2-12) is *nearly* 0, for some constants  $c_1, \dots, c_k$  the computer may produce incorrect inverses due to extreme errors in rounding. It is always good to check the products  $AA^{-1}$  and  $A^{-1}A$  for equality with  $I$  when  $A^{-1}$  is produced by a computer package. (See Exercise 2.10.)

Diagonal matrices have inverses that are easy to compute. For example,

$$\begin{bmatrix} a_{11} & 0 & 0 & 0 & 0 \\ 0 & a_{22} & 0 & 0 & 0 \\ 0 & 0 & a_{33} & 0 & 0 \\ 0 & 0 & 0 & a_{44} & 0 \\ 0 & 0 & 0 & 0 & a_{55} \end{bmatrix} \text{ has inverse } \begin{bmatrix} \frac{1}{a_{11}} & 0 & 0 & 0 & 0 \\ 0 & \frac{1}{a_{22}} & 0 & 0 & 0 \\ 0 & 0 & \frac{1}{a_{33}} & 0 & 0 \\ 0 & 0 & 0 & \frac{1}{a_{44}} & 0 \\ 0 & 0 & 0 & 0 & \frac{1}{a_{55}} \end{bmatrix}$$

if all the  $a_{ii} \neq 0$ .

Another special class of square matrices with which we shall become familiar are the *orthogonal* matrices, characterized by

$$QQ' = Q'Q = I \quad \text{or} \quad Q' = Q^{-1} \quad (2-13)$$

The name derives from the property that if  $Q$  has  $i$ th row  $q_i'$ , then  $QQ' = I$  implies that  $q_i'q_i = 1$  and  $q_i'q_j = 0$  for  $i \neq j$  so the rows have unit length and are mutually perpendicular (orthogonal). According to the condition  $Q'Q = I$ , the columns have the same property.

We conclude our brief introduction to the elements of matrix algebra by

introducing a concept fundamental to multivariate statistical analysis. A square matrix  $A$  is said to have an *eigenvalue*  $\lambda$ , with corresponding *eigenvector*  $x \neq 0$ , if

$$Ax = \lambda x \quad (2-14)$$

Ordinarily, we normalize  $x$  so that it has length one so  $1 = x'x$ . It is convenient to denote normalized eigenvectors by  $e$ , and we do so in the sequel. Sparing you the details of the derivation (see [1]), we state the following basic result.

Let  $A$  be a  $(k \times k)$  square symmetric matrix. Then  $A$  has  $k$  pairs of eigenvalues and eigenvectors:

$$\lambda_1, e_1 \quad \lambda_2, e_2 \quad \dots \quad \lambda_k, e_k \quad (2-15)$$

The eigenvectors can be chosen to satisfy  $1 = e_i'e_i = \dots = e_i'e_i$  and be mutually perpendicular. The eigenvectors are unique unless two or more eigenvalues are equal.

#### Example 2.8

Let

$$A = \begin{bmatrix} 1 & -5 \\ -5 & 1 \end{bmatrix}$$

Then, since

$$\begin{bmatrix} 1 & -5 \\ -5 & 1 \end{bmatrix} \begin{bmatrix} \frac{1}{\sqrt{2}} \\ -\frac{1}{\sqrt{2}} \end{bmatrix} = 6 \begin{bmatrix} \frac{1}{\sqrt{2}} \\ -\frac{1}{\sqrt{2}} \end{bmatrix}$$

$\lambda_1 = 6$  is an eigenvalue and

$$e_1 = \begin{bmatrix} \frac{1}{\sqrt{2}} \\ -\frac{1}{\sqrt{2}} \end{bmatrix}$$

is its corresponding normalized eigenvector. You may wish to show that a second eigenvalue, eigenvector pair is:  $\lambda_2 = -4$ ,  $e_2' = [1/\sqrt{2}, 1/\sqrt{2}]$ . ■

A method for calculating the  $\lambda$ 's and  $e$ 's is described in Supplement 2A. It is instructive to do a few sample calculations to understand the technique. We usually rely on a computer when the dimension of the square matrix is greater than two or three.



## 2.3 POSITIVE DEFINITE MATRICES

The study of the variation and interrelationships in multivariate data is often based upon distances and the assumption that the data are multivariate normally distributed. Squared distances (see Chapter 1) and the multivariate normal density can be expressed in terms of matrix products called *quadratic forms* (see Chapter 4). Consequently, it should not be surprising that quadratic forms play a central role in multivariate analysis. In this section we consider quadratic forms that are always nonnegative and the associated *positive definite* matrices.

Results involving quadratic forms and symmetric matrices are, in many cases, a direct consequence of an expansion for symmetric matrices known as the *spectral decomposition*. The spectral decomposition of a  $k \times k$  symmetric matrix  $A$  is given by<sup>1</sup>

$$A = \lambda_1 \mathbf{e}_1 \mathbf{e}_1' + \lambda_2 \mathbf{e}_2 \mathbf{e}_2' + \cdots + \lambda_k \mathbf{e}_k \mathbf{e}_k' \quad (2.16)$$

where  $\lambda_1, \lambda_2, \dots, \lambda_k$  are the eigenvalues of  $A$  and  $\mathbf{e}_1, \mathbf{e}_2, \dots, \mathbf{e}_k$  are the associated normalized eigenvectors. Thus  $\mathbf{e}_i' \mathbf{e}_j = 1$ ,  $i = 1, 2, \dots, k$  and  $\mathbf{e}_i' \mathbf{e}_j = 0$ ,  $i \neq j$ .

### Example 2.9

Consider the symmetric matrix

$$A = \begin{bmatrix} 13 & -4 & 2 \\ -4 & 13 & -2 \\ 2 & -2 & 10 \end{bmatrix}$$

The eigenvalues obtained from the characteristic equation  $|A - \lambda I| = 0$  are  $\lambda_1 = 9$ ,  $\lambda_2 = 9$ , and  $\lambda_3 = 18$  (Definition 2A.30). The corresponding eigenvectors  $\mathbf{e}_1, \mathbf{e}_2$ , and  $\mathbf{e}_3$  are the (normalized) solutions of the equations  $A\mathbf{e}_i = \lambda_i \mathbf{e}_i$  for  $i = 1, 2, 3$ . Thus  $A\mathbf{e}_1 = \lambda_1 \mathbf{e}_1$  gives

$$\begin{bmatrix} 13 & -4 & 2 \\ -4 & 13 & -2 \\ 2 & -2 & 10 \end{bmatrix} \begin{bmatrix} e_{11} \\ e_{21} \\ e_{31} \end{bmatrix} = 9 \begin{bmatrix} e_{11} \\ e_{21} \\ e_{31} \end{bmatrix}$$

or

$$\begin{aligned} 13e_{11} - 4e_{21} + 2e_{31} &= 9e_{11} \\ -4e_{11} + 13e_{21} - 2e_{31} &= 9e_{21} \\ 2e_{11} - 2e_{21} + 10e_{31} &= 9e_{31} \end{aligned}$$

Moving the terms on the right of the equals sign to the left yields three homogeneous equations in three unknowns, but two of the equations are redundant. Selecting one of the equations and arbitrarily setting  $e_{11} = 1$

<sup>1</sup> A proof of Equation (2.16) is beyond the scope of this book. The interested reader will find a proof in [5], Chapter 9.

and  $e_{21} = 1$ , we find  $e_{31} = 0$ . Consequently, the normalized eigenvector is  $\mathbf{e}_1' = [1/\sqrt{1^2 + 1^2 + 0^2}, 1/\sqrt{1^2 + 1^2 + 0^2}, 0/\sqrt{1^2 + 1^2 + 0^2}] = [1/\sqrt{2}, 1/\sqrt{2}, 0]$  since the sum of the squares of its elements is unity. You may verify that  $\mathbf{e}_2' = [1/\sqrt{18}, -1/\sqrt{18}, -4/\sqrt{18}]$  is also an eigenvector for  $9 = \lambda_2$ , and  $\mathbf{e}_3' = [2/3, -2/3, 1/3]$  is the normalized eigenvector corresponding to the eigenvalue  $\lambda_3 = 18$ . Moreover,  $\mathbf{e}_i' \mathbf{e}_j = 0$ ,  $i \neq j$ .

The spectral decomposition of  $A$  is then

$$A = \lambda_1 \mathbf{e}_1 \mathbf{e}_1' + \lambda_2 \mathbf{e}_2 \mathbf{e}_2' + \lambda_3 \mathbf{e}_3 \mathbf{e}_3'$$

or

$$\begin{bmatrix} 13 & -4 & 2 \\ -4 & 13 & -2 \\ 2 & -2 & 10 \end{bmatrix} = 9 \begin{bmatrix} \frac{1}{\sqrt{2}} & \frac{1}{\sqrt{2}} & 0 \\ \frac{1}{\sqrt{2}} & \frac{1}{\sqrt{2}} & 0 \\ 0 & 0 & 0 \end{bmatrix} + 9 \begin{bmatrix} \frac{1}{\sqrt{18}} & -\frac{1}{\sqrt{18}} & -\frac{4}{\sqrt{18}} \\ -\frac{1}{\sqrt{18}} & \frac{1}{\sqrt{18}} & \frac{4}{\sqrt{18}} \\ \frac{2}{\sqrt{18}} & -\frac{2}{\sqrt{18}} & \frac{1}{\sqrt{18}} \end{bmatrix} + 18 \begin{bmatrix} \frac{2}{3} & -\frac{2}{3} & \frac{1}{3} \\ -\frac{2}{3} & \frac{2}{3} & -\frac{1}{3} \\ \frac{1}{3} & -\frac{1}{3} & \frac{2}{3} \end{bmatrix}$$

$$= 9 \begin{bmatrix} \frac{1}{2} & \frac{1}{2} & 0 \\ \frac{1}{2} & \frac{1}{2} & 0 \\ 0 & 0 & 0 \end{bmatrix} + 9 \begin{bmatrix} \frac{1}{18} & -\frac{1}{18} & -\frac{4}{18} \\ -\frac{1}{18} & \frac{1}{18} & \frac{4}{18} \\ \frac{2}{18} & -\frac{2}{18} & \frac{1}{18} \end{bmatrix} + 18 \begin{bmatrix} \frac{4}{9} & -\frac{4}{9} & \frac{2}{9} \\ -\frac{4}{9} & \frac{4}{9} & -\frac{2}{9} \\ \frac{2}{9} & -\frac{2}{9} & \frac{4}{9} \end{bmatrix}$$

as you may readily verify. ■

The spectral decomposition is an important analytical tool. With it we are very easily able to demonstrate certain statistical results. The first of these is a matrix explanation of distance, which we now develop.



When a  $(k \times k)$  symmetric matrix  $A$  is such that

$$0 \leq x'Ax \quad (2-17)$$

for all  $x' = [x_1, x_2, \dots, x_k]$ ,  $A$  is said to be *nonnegative definite*. If equality holds in (2-17) only for the vector  $x' = [0, 0, \dots, 0]$ , then  $A$  is said to be *positive definite*. In other words,  $A$  is positive definite if

$$0 < x'Ax \quad (2-18)$$

for all vectors  $x \neq 0$ . Because  $x'Ax$  has only squared terms,  $x_1^2$ , and product terms,  $x_1x_2$ , it is called a *quadratic form*.

#### Example 2.10

Show that the following quadratic form is positive definite.

$$3x_1^2 + 2x_2^2 - 2\sqrt{2}x_1x_2$$

To illustrate the general approach, we first write the quadratic form in matrix notation as

$$\begin{bmatrix} x_1 & x_2 \end{bmatrix} \begin{bmatrix} 3 & -\sqrt{2} \\ -\sqrt{2} & 2 \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \end{bmatrix} = x'Ax$$

By Definition 2A.30, the eigenvalues of  $A$  are the solutions of the equation  $|A - \lambda I| = 0$ , or  $(3 - \lambda)(2 - \lambda) - 2 = 0$ . The solutions are  $\lambda_1 = 4$  and  $\lambda_2 = 1$ . Using the spectral decomposition in (2-16), we can write

$$\begin{aligned} A &= \lambda_1 e_1 e_1' + \lambda_2 e_2 e_2' \\ &= 4e_1 e_1' + e_2 e_2' \end{aligned}$$

where  $e_1$  and  $e_2$  are the normalized and orthogonal eigenvectors associated with the eigenvalues  $\lambda_1 = 4$  and  $\lambda_2 = 1$ , respectively. Because 4 and 1 are scalars, premultiplication and postmultiplication of  $A$  by  $x'$  and  $x$ , respectively, where  $x' = [x_1, x_2]$  is any *nonzero* vector, give

$$\begin{aligned} x'A &= 4x'e_1 e_1' + x'e_2 e_2' \\ &= 4y_1^2 + y_2^2 \geq 0 \end{aligned}$$

with

$$y_1 = x'e_1 = e_1'x \quad \text{and} \quad y_2 = x'e_2 = e_2'x$$

We now show that  $y_1$  and  $y_2$  are not both zero and, consequently, that  $x'Ax = 4y_1^2 + y_2^2 > 0$ , or  $A$  is *positive definite*.

From the definitions of  $y_1$  and  $y_2$  we have

$$\begin{bmatrix} y_1 \\ y_2 \end{bmatrix} = \begin{bmatrix} e_1' \\ e_2' \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \end{bmatrix}$$

or

$$\begin{bmatrix} y_1 \\ y_2 \end{bmatrix} = E'x$$

Now  $E$  is an orthogonal matrix and hence has inverse  $E'$ . Thus  $x = E'y$ . But  $x$  is a nonzero vector and  $0 \neq x = E'y$  implies  $y \neq 0$ .

Using the spectral decomposition it is rather easy to show that a  $(k \times k)$  symmetric matrix  $A$  is a positive definite matrix if and only if every eigenvalue of  $A$  is positive (see Exercise 2.17). It is a nonnegative definite matrix if and only if all of its eigenvalues are greater than or equal to zero.

Assume for the moment that the  $p$  elements  $x_1, x_2, \dots, x_p$  of a vector  $x$  are realizations of  $p$  random variables  $X_1, X_2, \dots, X_p$ . As we pointed out in Chapter 1, we can regard these elements as the coordinates of a point in  $p$ -dimensional space and the "distance" of the point  $[x_1, x_2, \dots, x_p]$  to the origin can, and in this case should, be interpreted in terms of standard deviation units. In this way we can account for the inherent uncertainty (variability) in the observations. Points with the same associated "uncertainty" are regarded as the same distance from the origin.

If we use the distance formula introduced in Chapter 1 [see Equation (1-22)], the distance from the origin satisfies the general formula

$$(\text{distance})^2 = a_{11}x_1^2 + a_{22}x_2^2 + \dots + a_{pp}x_p^2$$

provided  $(\text{distance})^2 > 0$  for all  $[x_1, x_2, \dots, x_p] \neq [0, 0, \dots, 0]$ . Setting  $a_{ij} = a_{ji}$ ,  $i \neq j$ ,  $i, j = 1, 2, \dots, p$ , we have

$$0 < (\text{distance})^2 = \begin{bmatrix} x_1 & x_2 & \dots & x_p \end{bmatrix} \begin{bmatrix} a_{11} & a_{12} & \dots & a_{1p} \\ a_{21} & a_{22} & \dots & a_{2p} \\ \vdots & \vdots & \ddots & \vdots \\ a_{p1} & a_{p2} & \dots & a_{pp} \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_p \end{bmatrix}$$

or

$$0 < (\text{distance})^2 = x'Ax \quad \text{for } x \neq 0 \quad (2-19)$$

From (2-19) we see that the  $(p \times p)$  symmetric matrix  $A$  is positive definite. In summary, distance is determined from a positive definite quadratic form  $x'Ax$ . Conversely, a positive definite quadratic form can be interpreted as a squared distance.

*Comment.* Let the squared distance from the point  $x' = [x_1, x_2, \dots, x_p]$  to the origin be given by  $x'Ax$  where  $A$  is a  $(p \times p)$  symmetric and positive definite matrix. Then the squared distance from  $x$  to an arbitrary fixed point  $\mu' = [\mu_1, \mu_2, \dots, \mu_p]$  is given by the general expression  $(x - \mu)'A(x - \mu)$ .

Expressing distance as the square root of a positive definite quadratic form allows us to give a geometrical interpretation based on the eigenvalues and eigenvectors of the matrix  $A$ . For example, suppose  $p = 2$ . The points  $x' = [x_1, x_2]$  of constant distance  $c$  from the origin satisfy

$$x'Ax = a_{11}x_1^2 + a_{22}x_2^2 + 2a_{12}x_1x_2 = c^2$$

By the spectral decomposition, as in Example 2.10,

$$A = \lambda_1 e_1 e_1' + \lambda_2 e_2 e_2' \quad \text{so} \quad x'Ax = \lambda_1 (x'e_1)^2 + \lambda_2 (x'e_2)^2$$



Now  $c^2 = \lambda_1^2 + \lambda_2^2$  is an ellipse in  $y_1 = x'e_1$  and  $y_2 = x'e_2$  because  $\lambda_1, \lambda_2 > 0$  when  $A$  is positive definite (see Exercise 2.17). We easily verify that  $x = c\lambda_1^{-1/2}e_1$  satisfies  $x'Ax = \lambda_1(c\lambda_1^{-1/2}e_1)'e_1 = c^2$  and  $x = c\lambda_2^{-1/2}e_2$  gives the appropriate distance in the  $e_2$  direction. Thus the points at distance  $c$  lie on an ellipse whose axes are given by the eigenvectors of  $A$  with lengths proportional to the reciprocals of the square roots of the eigenvalues. The constant of proportionality is  $c$ . The situation is illustrated in Figure 2.6.

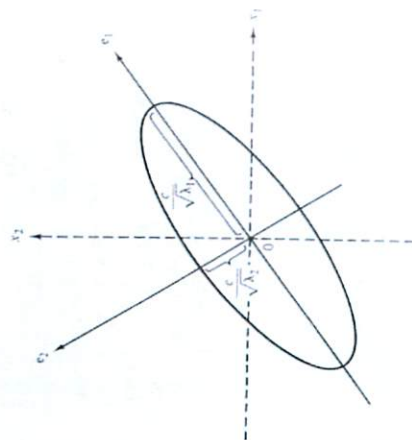


Figure 2.6 Points at a constant distance  $c$  from the origin ( $p = 2, 1 \leq \lambda_i < \lambda_j$ ).

If  $p > 2$ , the points  $x' = [x_1, x_2, \dots, x_p]$  at a constant distance  $c = \sqrt{x'Ax}$  from the origin lie on hyperellipsoids  $c^2 = \lambda_1(x'e_1)^2 + \dots + \lambda_p(x'e_p)^2$ , whose axes are given by the eigenvectors of  $A$ . The half-length in the direction  $e_i$  is equal to  $c/\sqrt{\lambda_i}$ ,  $i = 1, 2, \dots, p$ , where  $\lambda_1, \lambda_2, \dots, \lambda_p$  are the eigenvalues of  $A$ .

## 2.4 A SQUARE-ROOT MATRIX

The spectral decomposition allows us to express the inverse of a square matrix in terms of its eigenvalues and eigenvectors, and this leads to a useful square-root matrix.

Let  $A$  be a  $k \times k$  positive definite matrix with the spectral decomposition  $A = \sum_{i=1}^k \lambda_i e_i e_i'$ . Let the normalized eigenvectors be the columns of another matrix  $P = [e_1, e_2, \dots, e_k]$ . Then

$$A = \sum_{i=1}^k \lambda_i e_i e_i' = P \Lambda P' \quad (2.20)$$

50 Chap. 2 Matrix Algebra and Random Vectors

where  $PP' = P'P = I$  and  $\Lambda$  is the diagonal matrix

$$\Lambda = \begin{bmatrix} \lambda_1 & 0 & \dots & 0 \\ 0 & \lambda_2 & \dots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \dots & \lambda_k \end{bmatrix} \quad \text{with } \lambda_i > 0$$

Thus

$$A^{-1} = P\Lambda^{-1}P' = \sum_{i=1}^k \frac{1}{\lambda_i} e_i e_i' \quad (2.21)$$

since  $(P\Lambda^{-1}P')(P\Lambda P') = P\Lambda P'\Lambda^{-1}P' = PP' = I$ .

Next let  $\Lambda^{1/2}$  denote the diagonal matrix with  $\sqrt{\lambda_i}$  as the  $i$ th diagonal element. The matrix  $\sum_{i=1}^k \sqrt{\lambda_i} e_i e_i' = P\Lambda^{1/2}P'$  is called the square root of  $A$  and is denoted by  $A^{1/2}$ .

The square root matrix

$$A^{1/2} = \sum_{i=1}^k \sqrt{\lambda_i} e_i e_i' = P\Lambda^{1/2}P' \quad (2.22)$$

has the following properties:

1.  $(A^{1/2})' = A^{1/2}$  (that is,  $A^{1/2}$  is symmetric).

2.  $A^{1/2}A^{1/2} = A$

3.  $(A^{1/2})^{-1} = \sum_{i=1}^k \frac{1}{\sqrt{\lambda_i}} e_i e_i' = P\Lambda^{-1/2}P'$ , where  $\Lambda^{-1/2}$  is a diagonal matrix with  $1/\sqrt{\lambda_i}$  as the  $i$ th diagonal element.

4.  $A^{1/2}A^{-1/2} = A^{-1/2}A^{1/2} = I$  and  $A^{-1/2}A^{-1/2} = A^{-1}$  where  $A^{-1/2} = (A^{1/2})^{-1}$ .

## 2.5 RANDOM VECTORS AND MATRICES

A random vector is a vector whose elements are random variables. Similarly, a random matrix is a matrix whose elements are random variables. The expected value of a random matrix (or vector) is the matrix consisting of the expected values of each of its elements. Specifically, let  $X = [X_{ij}]$  be a  $(p \times n)$  random matrix. The expected value of  $X$ , denoted by  $E(X)$ , is the  $(p \times n)$  matrix of numbers (if they exist)



$$E(\mathbf{X}) = \begin{bmatrix} E(X_{11}) & E(X_{12}) & \cdots & E(X_{1n}) \\ E(X_{21}) & E(X_{22}) & \cdots & E(X_{2n}) \\ \vdots & \vdots & \ddots & \vdots \\ E(X_{p1}) & E(X_{p2}) & \cdots & E(X_{pn}) \end{bmatrix} \quad (2-23)$$

where, for each element of the matrix<sup>2</sup>

$$E(X_{ij}) = \begin{cases} \int_{-\infty}^{\infty} x_{ij} f_{ij}(x_{ij}) dx_{ij} & \text{if } X_{ij} \text{ is a continuous random variable} \\ & \text{with probability density function } f_{ij}(x_{ij}) \\ \sum_{\text{all } x_{ij}} x_{ij} p_{ij}(x_{ij}) & \text{if } X_{ij} \text{ is a discrete random variable} \\ & \text{with probability function } p_{ij}(x_{ij}) \end{cases}$$

### Example 2.11

Suppose  $p = 2$  and  $n = 1$  and consider the random vector  $\mathbf{X}' = [X_1, X_2]$ . Let the discrete random variable  $X_1$  have the following probability function:

|            |    |    |    |
|------------|----|----|----|
| $x_1$      | -1 | 0  | .1 |
| $p_1(x_1)$ | .3 | .3 | .4 |

Then  $E(X_1) = \sum_{\text{all } x_1} x_1 p_1(x_1) = (-1)(.3) + (0)(.3) + (1)(.4) = .1$ .

Similarly, let the discrete random variable  $X_2$  have the probability function

|            |    |    |
|------------|----|----|
| $x_2$      | 0  | 1  |
| $p_2(x_2)$ | .8 | .2 |

Then  $E(X_2) = \sum_{\text{all } x_2} x_2 p_2(x_2) = (0)(.8) + (1)(.2) = .2$ .

Thus

$$E(\mathbf{X}) = \begin{bmatrix} E(X_1) \\ E(X_2) \end{bmatrix} = \begin{bmatrix} .1 \\ .2 \end{bmatrix}$$

Two results involving the expectation of sums and products of matrices follow directly from the definition of the expected value of a random matrix and the univariate properties of expectation,  $E(X_1 + Y_1) = E(X_1) + E(Y_1)$  and  $E(cX_1) = cE(X_1)$ .

Let  $\mathbf{X}$  and  $\mathbf{Y}$  be random matrices of the same dimension and let  $\mathbf{A}$  and  $\mathbf{B}$  be conformable matrices of constants. Then (see Exercise 2.35)

$$\begin{aligned} E(\mathbf{X} + \mathbf{Y}) &= E(\mathbf{X}) + E(\mathbf{Y}) \\ E(\mathbf{AXB}) &= \mathbf{AE(X)B} \end{aligned} \quad (2-24)$$

<sup>2</sup> If you are unfamiliar with calculus, you should concentrate on the interpretation of the expected value and, eventually, variance. Our development is based primarily on the properties of expectation rather than its particular evaluation for continuous or discrete random variables.

## 2.6 MEAN VECTORS AND COVARIANCE MATRICES

Suppose  $\mathbf{X} = \{X_i\}$  is a  $(p \times 1)$  random matrix, that is, a random vector. Each element of  $\mathbf{X}$  is a random variable with its own marginal probability distribution (see Example 2.11). The marginal means,  $\mu_i$ , and variances,  $\sigma_i^2$ , are defined as  $\mu_i = E(X_i)$  and  $\sigma_i^2 = E(X_i - \mu_i)^2$ ,  $i = 1, 2, \dots, p$ , respectively. Specifically,

$$\mu_i = \begin{cases} \int_{-\infty}^{\infty} x_i f_i(x_i) dx_i & \text{if } X_i \text{ is a continuous random variable with} \\ & \text{probability density function } f_i(x_i) \\ \sum_{\text{all } x_i} x_i p_i(x_i) & \text{if } X_i \text{ is a discrete random variable with} \\ & \text{probability function } p_i(x_i) \end{cases}$$

$$\sigma_i^2 = \begin{cases} \int_{-\infty}^{\infty} (x_i - \mu_i)^2 f_i(x_i) dx_i & \text{if } X_i \text{ is a continuous random variable with} \\ & \text{probability density function } f_i(x_i) \\ \sum_{\text{all } x_i} (x_i - \mu_i)^2 p_i(x_i) & \text{if } X_i \text{ is a discrete random variable} \\ & \text{with probability function } p_i(x_i) \end{cases} \quad (2-25)$$

It will be convenient in later sections to denote the marginal variances by  $\sigma_{ii}$  rather than the more traditional  $\sigma_i^2$  and, consequently, we shall adopt this notation.

The behavior of any pair of random variables, such as  $X_i$  and  $X_k$ , is described by their joint probability function and a measure of the linear association between them is provided by the covariance  $\sigma_{ik}$ , where

$$\sigma_{ik} = \begin{cases} E(X_i - \mu_i)(X_k - \mu_k) \\ \int_{-\infty}^{\infty} \int_{-\infty}^{\infty} (x_i - \mu_i)(x_k - \mu_k) f_{ik}(x_i, x_k) dx_i dx_k & \text{if } X_i, X_k \text{ are} \\ & \text{continuous random} \\ & \text{variables with the} \\ & \text{joint density} \\ & \text{function } f_{ik}(x_i, x_k) \\ \sum_{\text{all } x_i, \text{ all } x_k} (x_i - \mu_i)(x_k - \mu_k) p_{ik}(x_i, x_k) & \text{if } X_i, X_k \text{ are} \\ & \text{discrete random} \\ & \text{variables with the} \\ & \text{joint probability} \\ & \text{function } p_{ik}(x_i, x_k) \end{cases} \quad (2-26)$$

and  $\mu_i$  and  $\mu_k$ ,  $i, k = 1, 2, \dots, p$ , are the marginal means. When  $i = k$ , the covariance becomes the marginal variance.

More generally, the collective behavior of the  $p$  random variables  $X_1, X_2, \dots, X_p$ , or, equivalently, the random vector  $\mathbf{X}' = [X_1, X_2, \dots, X_p]$  is described by a joint probability density function  $f(x_1, x_2, \dots, x_p) = f(\mathbf{x})$ . As we have already noted in this book,  $f(\mathbf{x})$  will often be the multivariate normal density function (see Chapter 4).

If the joint probability  $P[X_i \leq x_i \text{ and } X_k \leq x_k]$  can be written as the product of the corresponding marginal probabilities so that

$$P[X_i \leq x_i \text{ and } X_k \leq x_k] = P[X_i \leq x_i] P[X_k \leq x_k] \quad (2-27)$$



for all pairs of values  $x_i, x_k$ , then  $X_i$  and  $X_k$  are said to be *statistically independent*. When  $X_i$  and  $X_k$  are continuous random variables with joint density  $f_{ik}(x_i, x_k)$  and marginal densities  $f_i(x_i)$  and  $f_k(x_k)$ , the independence condition becomes

$$f_{ik}(x_i, x_k) = f_i(x_i)f_k(x_k)$$

for all pairs  $(x_i, x_k)$ .

The  $p$  continuous random variables  $X_1, X_2, \dots, X_p$  are *mutually statistically independent* if their joint density factors as

$$f_{12\dots p}(x_1, x_2, \dots, x_p) = f_1(x_1)f_2(x_2) \cdots f_p(x_p) \quad (2-28)$$

for all  $p$ -tuples  $(x_1, x_2, \dots, x_p)$ .

Statistical independence has an important implication for covariance. The factorization in (2-28) implies  $\text{Cov}(X_i, X_k) = 0$ . Thus

$$\text{Cov}(X_i, X_k) = 0 \quad \text{if } X_i \text{ and } X_k \text{ are independent.} \quad (2-29)$$

The converse of (2-29) is not true in general. There are situations where  $\text{Cov}(X_i, X_k) = 0$  and  $X_i$  and  $X_k$  are not independent (see [2]).

The means and covariances of the  $(p \times 1)$  random vector  $\mathbf{X}$  can be set out as matrices. The expected value of each element is contained in the vector of means  $\boldsymbol{\mu} = E(\mathbf{X})$  and the  $p$  variances  $\sigma_{ii}$  and the  $p(p-1)/2$  distinct covariances  $\sigma_{ik}$  ( $i < k$ ) are contained in the symmetric variance-covariance matrix  $\boldsymbol{\Sigma} = E(\mathbf{X} - \boldsymbol{\mu})(\mathbf{X} - \boldsymbol{\mu})'$ . Specifically,

$$E(\mathbf{X}) = \begin{bmatrix} E(X_1) \\ E(X_2) \\ \vdots \\ E(X_p) \end{bmatrix} = \begin{bmatrix} \mu_1 \\ \mu_2 \\ \vdots \\ \mu_p \end{bmatrix} = \boldsymbol{\mu} \quad (2-30)$$

and

$$\begin{aligned} \boldsymbol{\Sigma} &= E(\mathbf{X} - \boldsymbol{\mu})(\mathbf{X} - \boldsymbol{\mu})' = E \left( \begin{bmatrix} X_1 - \mu_1 \\ X_2 - \mu_2 \\ \vdots \\ X_p - \mu_p \end{bmatrix} [X_1 - \mu_1, X_2 - \mu_2, \dots, X_p - \mu_p] \right) \\ &= E \begin{bmatrix} (X_1 - \mu_1)^2 & (X_1 - \mu_1)(X_2 - \mu_2) & \cdots & (X_1 - \mu_1)(X_p - \mu_p) \\ (X_2 - \mu_2)(X_1 - \mu_1) & (X_2 - \mu_2)^2 & \cdots & (X_2 - \mu_2)(X_p - \mu_p) \\ \vdots & \vdots & \ddots & \vdots \\ (X_p - \mu_p)(X_1 - \mu_1) & (X_p - \mu_p)(X_2 - \mu_2) & \cdots & (X_p - \mu_p)^2 \end{bmatrix} \\ &= \begin{bmatrix} E(X_1 - \mu_1)^2 & E(X_1 - \mu_1)(X_2 - \mu_2) & \cdots & E(X_1 - \mu_1)(X_p - \mu_p) \\ E(X_2 - \mu_2)(X_1 - \mu_1) & E(X_2 - \mu_2)^2 & \cdots & E(X_2 - \mu_2)(X_p - \mu_p) \\ \vdots & \vdots & \ddots & \vdots \\ E(X_p - \mu_p)(X_1 - \mu_1) & E(X_p - \mu_p)(X_2 - \mu_2) & \cdots & E(X_p - \mu_p)^2 \end{bmatrix} \end{aligned}$$

or

$$\boldsymbol{\Sigma} = \text{Cov}(\mathbf{X}) = \begin{bmatrix} \sigma_{11} & \sigma_{12} & \cdots & \sigma_{1p} \\ \sigma_{21} & \sigma_{22} & \cdots & \sigma_{2p} \\ \vdots & \vdots & \ddots & \vdots \\ \sigma_{p1} & \sigma_{p2} & \cdots & \sigma_{pp} \end{bmatrix} \quad (2-31)$$

### Example 2.12

Find the covariance matrix for the two random variables  $X_1$  and  $X_2$  introduced in Example 2.11 when their joint probability function,  $p_{12}(x_1, x_2)$  is represented by the entries in the body of the following table.

| $x_1 \backslash x_2$ | 0   | 1   | $p_1(x_1)$ |
|----------------------|-----|-----|------------|
| -1                   | .24 | .06 | .3         |
| 0                    | .16 | .14 | .3         |
| 1                    | .40 | .00 | .4         |
| $p_2(x_2)$           | .8  | .2  | 1          |

We have already shown that  $\mu_1 = E(X_1) = .1$  and  $\mu_2 = E(X_2) = .2$  (see Example 2.11). In addition

$$\sigma_{11} = E(X_1 - \mu_1)^2 = \sum_{\text{all } x_1} (x_1 - .1)^2 p_1(x_1)$$

$$= (-1 - .1)^2(.3) + (0 - .1)^2(.3) + (1 - .1)^2(.4) = .69$$

$$\sigma_{22} = E(X_2 - \mu_2)^2 = \sum_{\text{all } x_2} (x_2 - .2)^2 p_2(x_2)$$

$$= (0 - .2)^2(.8) + (1 - .2)^2(.2)$$

$$= .16$$

$$\sigma_{12} = E(X_1 - \mu_1)(X_2 - \mu_2) = \sum_{\text{all pairs } (x_1, x_2)} (x_1 - .1)(x_2 - .2)p_{12}(x_1, x_2)$$

$$= (-1 - .1)(0 - .2)(.24) + (-1 - .1)(1 - .2)(.06) + \cdots + (1 - .1)(1 - .2)(.00) = -.08$$

$$\sigma_{21} = E(X_2 - \mu_2)(X_1 - \mu_1) = E(X_1 - \mu_1)(X_2 - \mu_2) = \sigma_{12} = -.08$$

Consequently, with  $\mathbf{X}' = [X_1, X_2]$ ,

$$\boldsymbol{\mu} = E(\mathbf{X}) = \begin{bmatrix} E(X_1) \\ E(X_2) \end{bmatrix} = \begin{bmatrix} \mu_1 \\ \mu_2 \end{bmatrix} = \begin{bmatrix} .1 \\ .2 \end{bmatrix}$$

and

$$\begin{aligned} \boldsymbol{\Sigma} &= E(\mathbf{X} - \boldsymbol{\mu})(\mathbf{X} - \boldsymbol{\mu})' = E \begin{bmatrix} (X_1 - \mu_1)^2 & (X_1 - \mu_1)(X_2 - \mu_2) \\ (X_2 - \mu_2)(X_1 - \mu_1) & (X_2 - \mu_2)^2 \end{bmatrix} \\ &= \begin{bmatrix} E(X_1 - \mu_1)^2 & E(X_1 - \mu_1)(X_2 - \mu_2) \\ E(X_2 - \mu_2)(X_1 - \mu_1) & E(X_2 - \mu_2)^2 \end{bmatrix} \\ &= \begin{bmatrix} \sigma_{11} & \sigma_{12} \\ \sigma_{21} & \sigma_{22} \end{bmatrix} = \begin{bmatrix} .69 & -.08 \\ -.08 & .16 \end{bmatrix} \end{aligned}$$



We note that the computation of means, variances, and covariances for discrete random variables involves summation (as in Examples 2.11 and 2.12), while analogous computations for continuous random variables involve integration.

Because  $\sigma_{ik} = E(X_i - \mu)(X_k - \mu) = \sigma_{ki}$ , it is convenient to write the matrix appearing in (2-31) as

$$\Sigma = E(X - \mu)(X - \mu)' = \begin{bmatrix} \sigma_{11} & \sigma_{12} & \cdots & \sigma_{1p} \\ \sigma_{12} & \sigma_{22} & \cdots & \sigma_{2p} \\ \vdots & \vdots & \ddots & \vdots \\ \sigma_{1p} & \sigma_{2p} & \cdots & \sigma_{pp} \end{bmatrix} \quad (2-32)$$

We shall refer to  $\mu$  and  $\Sigma$  as the *population mean* (vector) and *population variance-covariance* (matrix), respectively.

The multivariate normal distribution is completely specified once the mean vector  $\mu$  and variance-covariance matrix  $\Sigma$  are given (see Chapter 4), so it is not surprising that these quantities play an important role in many multivariate procedures.

It is frequently informative to separate the information contained in variances,  $\sigma_{ii}$ , from that contained in measures of association and, in particular, the measure of association known as the *population correlation coefficient*,  $\rho_{ik}$ . The correlation coefficient  $\rho_{ik}$  is defined in terms of the covariance  $\sigma_{ik}$  and variances  $\sigma_{ii}$  and  $\sigma_{kk}$  as

$$\rho_{ik} = \frac{\sigma_{ik}}{\sqrt{\sigma_{ii}\sigma_{kk}}} \quad (2-33)$$

The correlation coefficient measures the amount of linear association between the random variables  $X_i$  and  $X_k$ . (See, for example, [2].)

Let the population correlation matrix be the  $(p \times p)$  symmetric matrix  $\rho$  where

$$\rho = \begin{bmatrix} \frac{\sigma_{11}}{\sqrt{\sigma_{11}\sigma_{11}}} & \frac{\sigma_{12}}{\sqrt{\sigma_{11}\sigma_{22}}} & \cdots & \frac{\sigma_{1p}}{\sqrt{\sigma_{11}\sigma_{pp}}} \\ \frac{\sigma_{12}}{\sqrt{\sigma_{11}\sigma_{22}}} & \frac{\sigma_{22}}{\sqrt{\sigma_{22}\sigma_{22}}} & \cdots & \frac{\sigma_{2p}}{\sqrt{\sigma_{22}\sigma_{pp}}} \\ \vdots & \vdots & \ddots & \vdots \\ \frac{\sigma_{1p}}{\sqrt{\sigma_{11}\sigma_{pp}}} & \frac{\sigma_{2p}}{\sqrt{\sigma_{22}\sigma_{pp}}} & \cdots & \frac{\sigma_{pp}}{\sqrt{\sigma_{pp}\sigma_{pp}}} \end{bmatrix} = \begin{bmatrix} 1 & \rho_{12} & \cdots & \rho_{1p} \\ \rho_{12} & 1 & \cdots & \rho_{2p} \\ \vdots & \vdots & \ddots & \vdots \\ \rho_{1p} & \rho_{2p} & \cdots & 1 \end{bmatrix} \quad (2-34)$$

and let the  $(p \times p)$  standard deviation matrix,  $V^{1/2}$ , be

$$V^{1/2} = \begin{bmatrix} \sqrt{\sigma_{11}} & 0 & \cdots & 0 \\ 0 & \sqrt{\sigma_{22}} & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & \sqrt{\sigma_{pp}} \end{bmatrix} \quad (2-35)$$

Then it is easily verified (see Exercise 2.21) that

$$V^{1/2} \rho V^{1/2} = \Sigma \quad (2-36)$$

and

$$\rho = (V^{1/2})^{-1} \Sigma (V^{1/2})^{-1} \quad (2-37)$$

That is,  $\Sigma$  can be obtained from  $V^{1/2}$  and  $\rho$ , while  $\rho$  can be obtained from  $\Sigma$ . Moreover, the expression of these relationships in terms of matrix operations allows the calculations to be conveniently implemented on a computer.

### Example 2.13

Suppose

$$\Sigma = \begin{bmatrix} 4 & 1 & 2 \\ 1 & 9 & -3 \\ 2 & -3 & 25 \end{bmatrix} = \begin{bmatrix} \sigma_{11} & \sigma_{12} & \sigma_{13} \\ \sigma_{12} & \sigma_{22} & \sigma_{23} \\ \sigma_{13} & \sigma_{23} & \sigma_{33} \end{bmatrix}$$

Obtain  $V^{1/2}$  and  $\rho$ .

Here

$$V^{1/2} = \begin{bmatrix} \sqrt{\sigma_{11}} & 0 & 0 \\ 0 & \sqrt{\sigma_{22}} & 0 \\ 0 & 0 & \sqrt{\sigma_{33}} \end{bmatrix} = \begin{bmatrix} 2 & 0 & 0 \\ 0 & 3 & 0 \\ 0 & 0 & 5 \end{bmatrix}$$

and

$$(V^{1/2})^{-1} = \begin{bmatrix} \frac{1}{2} & 0 & 0 \\ 0 & \frac{1}{3} & 0 \\ 0 & 0 & \frac{1}{5} \end{bmatrix}$$

Consequently, using (2-37), the correlation matrix  $\rho$  is given by

$$(V^{1/2})^{-1} \Sigma (V^{1/2})^{-1} = \begin{bmatrix} \frac{1}{2} & 0 & 0 \\ 0 & \frac{1}{3} & 0 \\ 0 & 0 & \frac{1}{5} \end{bmatrix} \begin{bmatrix} 4 & 1 & 2 \\ 9 & -3 & 0 \\ 2 & -3 & 25 \end{bmatrix} \begin{bmatrix} \frac{1}{2} & 0 & 0 \\ 0 & \frac{1}{3} & 0 \\ 0 & 0 & \frac{1}{5} \end{bmatrix} = \begin{bmatrix} 1 & \frac{1}{6} & \frac{1}{5} \\ \frac{1}{6} & 1 & -\frac{1}{5} \\ \frac{1}{5} & -\frac{1}{5} & 1 \end{bmatrix}$$

### Partitioning the Covariance Matrix

Often, the characteristics measured on individual trials will fall naturally into two or more groups. As examples, consider measurements of variables representing consumption and income or variables representing personality traits and physical characteristics. One approach to handling these situations is to let the characteristics



defining the distinct groups be subsets of the total collection of characteristics. If the total collection is represented by a  $(p \times 1)$ -dimensional random vector  $X$ , the subsets can be regarded as components of  $X$  and can be sorted by partitioning  $X$ .

In general, we can partition the  $p$  characteristics contained in the  $p \times 1$  random vector  $X$  into, for instance, two groups of size  $q$  and  $p - q$ , respectively. For example, we can write

$$X = \begin{bmatrix} X_1 \\ \vdots \\ X_q \\ \vdots \\ X_{q+1} \\ \vdots \\ X_p \end{bmatrix} = \begin{bmatrix} X^{(1)} \\ X^{(2)} \end{bmatrix} \quad \text{and} \quad \mu = E(X) = \begin{bmatrix} \mu_1 \\ \vdots \\ \mu_q \\ \vdots \\ \mu_{q+1} \\ \vdots \\ \mu_p \end{bmatrix} = \begin{bmatrix} \mu^{(1)} \\ \mu^{(2)} \end{bmatrix} \quad (2.38)$$

From the definitions of transpose and matrix multiplication,

$$(X^{(1)} - \mu^{(1)})(X^{(2)} - \mu^{(2)})' = \begin{bmatrix} X_1 - \mu_1 \\ X_2 - \mu_2 \\ \vdots \\ X_q - \mu_q \end{bmatrix} \begin{bmatrix} X_{q+1} - \mu_{q+1} & X_{q+2} - \mu_{q+2} & \cdots & X_p - \mu_p \end{bmatrix} =$$

$$\begin{bmatrix} (X_1 - \mu_1)(X_{q+1} - \mu_{q+1}) & (X_1 - \mu_1)(X_{q+2} - \mu_{q+2}) & \cdots & (X_1 - \mu_1)(X_p - \mu_p) \\ (X_2 - \mu_2)(X_{q+1} - \mu_{q+1}) & (X_2 - \mu_2)(X_{q+2} - \mu_{q+2}) & \cdots & (X_2 - \mu_2)(X_p - \mu_p) \\ \vdots & \vdots & \ddots & \vdots \\ (X_q - \mu_q)(X_{q+1} - \mu_{q+1}) & (X_q - \mu_q)(X_{q+2} - \mu_{q+2}) & \cdots & (X_q - \mu_q)(X_p - \mu_p) \end{bmatrix}$$

Upon taking the expectation of the matrix  $(X^{(1)} - \mu^{(1)})(X^{(2)} - \mu^{(2)})'$ , we get

$$E[(X^{(1)} - \mu^{(1)})(X^{(2)} - \mu^{(2)})'] = \begin{bmatrix} \sigma_{1,q+1} & \sigma_{1,q+2} & \cdots & \sigma_{1p} \\ \sigma_{2,q+1} & \sigma_{2,q+2} & \cdots & \sigma_{2p} \\ \vdots & \vdots & \ddots & \vdots \\ \sigma_{q,q+1} & \sigma_{q,q+2} & \cdots & \sigma_{qp} \end{bmatrix} = \Sigma_{12} \quad (2.39)$$

which gives all the covariances,  $\sigma_{ij}$ ,  $i = 1, 2, \dots, q$ ,  $j = q + 1, q + 2, \dots, p$ , between a component of  $X^{(1)}$  and a component of  $X^{(2)}$ . You will note that the matrix  $\Sigma_{12}$  is not necessarily symmetric or even square.

Making use of the partitioning in Equation (2.38), it is easily demonstrated that

$$(X - \mu)(X - \mu)' = \begin{bmatrix} (X^{(1)} - \mu^{(1)})(X^{(1)} - \mu^{(1)})' & (X^{(1)} - \mu^{(1)})(X^{(2)} - \mu^{(2)})' \\ (X^{(2)} - \mu^{(2)})(X^{(1)} - \mu^{(1)})' & (X^{(2)} - \mu^{(2)})(X^{(2)} - \mu^{(2)})' \end{bmatrix}$$

and consequently

$$\Sigma_{(p \times p)} = E[(X - \mu)(X - \mu)'] = \begin{bmatrix} \Sigma_{11} & \Sigma_{12} \\ \Sigma_{21} & \Sigma_{22} \end{bmatrix} \quad (2.40)$$

$$= \begin{bmatrix} \sigma_{11} & \cdots & \sigma_{1q} & \cdots & \sigma_{1,q+1} & \cdots & \sigma_{1p} \\ \vdots & \ddots & \vdots & \ddots & \vdots & \ddots & \vdots \\ \sigma_{q1} & \cdots & \sigma_{qq} & \cdots & \sigma_{q,q+1} & \cdots & \sigma_{qp} \\ \vdots & \ddots & \vdots & \ddots & \vdots & \ddots & \vdots \\ \sigma_{q+1,1} & \cdots & \sigma_{q+1,q} & \cdots & \sigma_{q+1,q+1} & \cdots & \sigma_{q+1,p} \\ \vdots & \ddots & \vdots & \ddots & \vdots & \ddots & \vdots \\ \sigma_{p1} & \cdots & \sigma_{pq} & \cdots & \sigma_{p,q+1} & \cdots & \sigma_{pp} \end{bmatrix}$$

Note that  $\Sigma_{12} = \Sigma_{21}'$ . The covariance matrix of  $X^{(1)}$  is  $\Sigma_{11}$ , that of  $X^{(2)}$  is  $\Sigma_{22}$ , and that of elements from  $X^{(1)}$  and  $X^{(2)}$  is  $\Sigma_{12}$  (or  $\Sigma_{21}'$ ).

### The Mean Vector and Covariance Matrix for Linear Combinations of Random Variables

Recall that if a single random variable, such as  $X_1$ , is multiplied by a constant  $c$ , then

$$E(cX_1) = cE(X_1) = c\mu_1$$

and

$$\text{Var}(cX_1) = E[(cX_1 - c\mu_1)^2] = c^2 \text{Var}(X_1) = c^2 \sigma_{11}$$

If  $X_2$  is a second random variable and  $a$  and  $b$  are constants, then using additional properties of expectation

$$\begin{aligned} \text{Cov}(aX_1, bX_2) &= E[a(X_1 - \mu_1)](b(X_2 - \mu_2)) \\ &= abE[(X_1 - \mu_1)(X_2 - \mu_2)] \\ &= ab \text{Cov}(X_1, X_2) = ab\sigma_{12} \end{aligned}$$

Finally, for the linear combination  $aX_1 + bX_2$ , we have

$$\begin{aligned} E(aX_1 + bX_2) &= aE(X_1) + bE(X_2) = a\mu_1 + b\mu_2 \\ \text{Var}(aX_1 + bX_2) &= E[(aX_1 + bX_2) - (a\mu_1 + b\mu_2)]^2 \\ &= E[a(X_1 - \mu_1) + b(X_2 - \mu_2)]^2 \\ &= E[a^2(X_1 - \mu_1)^2 + b^2(X_2 - \mu_2)^2 + 2ab(X_1 - \mu_1)(X_2 - \mu_2)] \\ &= a^2 \text{Var}(X_1) + b^2 \text{Var}(X_2) + 2ab \text{Cov}(X_1, X_2) \\ &= a^2 \sigma_{11} + b^2 \sigma_{22} + 2ab\sigma_{12} \end{aligned} \quad (2.41)$$



With  $\mathbf{c}' = [a, b]$ ,  $aX_1 + bX_2$  can be written as

$$\begin{bmatrix} a & b \end{bmatrix} \begin{bmatrix} X_1 \\ X_2 \end{bmatrix} = \mathbf{c}'\mathbf{X}$$

Similarly,  $E(aX_1 + bX_2) = a\mu_1 + b\mu_2$  can be expressed as

$$\begin{bmatrix} a & b \end{bmatrix} \begin{bmatrix} \mu_1 \\ \mu_2 \end{bmatrix} = \mathbf{c}'\boldsymbol{\mu}$$

If we let

$$\boldsymbol{\Sigma} = \begin{bmatrix} \sigma_{11} & \sigma_{12} \\ \sigma_{12} & \sigma_{22} \end{bmatrix}$$

be the variance-covariance matrix of  $\mathbf{X}$ , Equation (2.41) becomes

$$\text{Var}(aX_1 + bX_2) = \text{Var}(\mathbf{c}'\mathbf{X}) = \mathbf{c}'\boldsymbol{\Sigma}\mathbf{c} \quad (2.42)$$

since

$$\mathbf{c}'\boldsymbol{\Sigma}\mathbf{c} = \begin{bmatrix} a & b \end{bmatrix} \begin{bmatrix} \sigma_{11} & \sigma_{12} \\ \sigma_{12} & \sigma_{22} \end{bmatrix} \begin{bmatrix} a \\ b \end{bmatrix} = a^2\sigma_{11} + 2ab\sigma_{12} + b^2\sigma_{22}$$

The results above can be extended to a linear combination of  $p$  random variables.

The linear combination  $\mathbf{c}'\mathbf{X} = c_1X_1 + \dots + c_pX_p$  has  
 mean  $= E(\mathbf{c}'\mathbf{X}) = \mathbf{c}'\boldsymbol{\mu}$   
 variance  $= \text{Var}(\mathbf{c}'\mathbf{X}) = \mathbf{c}'\boldsymbol{\Sigma}\mathbf{c}$   
 where  $\boldsymbol{\mu} = E(\mathbf{X})$  and  $\boldsymbol{\Sigma} = \text{Cov}(\mathbf{X})$ . (2.43)

In general consider the  $q$  linear combinations of the  $p$  random variables  $X_1, \dots, X_p$ ,

$$Z_1 = c_{11}X_1 + c_{12}X_2 + \dots + c_{1p}X_p$$

$$Z_2 = c_{21}X_1 + c_{22}X_2 + \dots + c_{2p}X_p$$

$$\vdots$$

$$Z_q = c_{q1}X_1 + c_{q2}X_2 + \dots + c_{qp}X_p$$

or

$$\mathbf{Z} = \begin{bmatrix} Z_1 \\ Z_2 \\ \vdots \\ Z_q \end{bmatrix} = \begin{bmatrix} c_{11} & c_{12} & \dots & c_{1p} \\ c_{21} & c_{22} & \dots & c_{2p} \\ \vdots & \vdots & \ddots & \vdots \\ c_{q1} & c_{q2} & \dots & c_{qp} \end{bmatrix} \begin{bmatrix} X_1 \\ X_2 \\ \vdots \\ X_p \end{bmatrix} = \mathbf{C}\mathbf{X} \quad (2.44)$$

$(q \times 1) \quad (q \times p) \quad (p \times 1)$

The linear combinations  $\mathbf{Z} = \mathbf{C}\mathbf{X}$  have

$$\begin{aligned} \boldsymbol{\mu}_Z &= E(\mathbf{Z}) = E(\mathbf{C}\mathbf{X}) = \mathbf{C}\boldsymbol{\mu}_X \\ \boldsymbol{\Sigma}_Z &= \text{Cov}(\mathbf{Z}) = \text{Cov}(\mathbf{C}\mathbf{X}) = \mathbf{C}\boldsymbol{\Sigma}_X\mathbf{C}' \end{aligned} \quad (2.45)$$

where  $\boldsymbol{\mu}_X$  and  $\boldsymbol{\Sigma}_X$  are the mean vector and variance-covariance matrix of  $\mathbf{X}$ , respectively. (See Exercise 2.26 for the computation of the off-diagonal terms in  $\mathbf{C}\boldsymbol{\Sigma}_X\mathbf{C}'$ .)

We shall rely heavily on the result in (2.45) in our discussions of principal components and factor analysis in Chapters 8 and 9.

#### Example 2.14

Let  $\mathbf{X}' = [X_1, X_2]$  be a random vector with mean vector  $\boldsymbol{\mu}'_X = [\mu_1, \mu_2]$  and variance-covariance matrix

$$\boldsymbol{\Sigma}_X = \begin{bmatrix} \sigma_{11} & \sigma_{12} \\ \sigma_{12} & \sigma_{22} \end{bmatrix}$$

Find the means and covariance matrix for the linear combinations

$$Z_1 = X_1 - X_2$$

$$Z_2 = X_1 + X_2$$

or

$$\mathbf{Z} = \begin{bmatrix} Z_1 \\ Z_2 \end{bmatrix} = \begin{bmatrix} 1 & -1 \\ 1 & 1 \end{bmatrix} \begin{bmatrix} X_1 \\ X_2 \end{bmatrix} = \mathbf{C}\mathbf{X}$$

in terms of  $\boldsymbol{\mu}_X$  and  $\boldsymbol{\Sigma}_X$ .

Here

$$\boldsymbol{\mu}_Z = E(\mathbf{Z}) = \mathbf{C}\boldsymbol{\mu}_X = \begin{bmatrix} 1 & -1 \\ 1 & 1 \end{bmatrix} \begin{bmatrix} \mu_1 \\ \mu_2 \end{bmatrix} = \begin{bmatrix} \mu_1 - \mu_2 \\ \mu_1 + \mu_2 \end{bmatrix}$$

and

$$\begin{aligned} \boldsymbol{\Sigma}_Z = \text{Cov}(\mathbf{Z}) &= \mathbf{C}\boldsymbol{\Sigma}_X\mathbf{C}' = \begin{bmatrix} 1 & -1 \\ 1 & 1 \end{bmatrix} \begin{bmatrix} \sigma_{11} & \sigma_{12} \\ \sigma_{12} & \sigma_{22} \end{bmatrix} \begin{bmatrix} 1 & 1 \\ -1 & 1 \end{bmatrix} \\ &= \begin{bmatrix} \sigma_{11} - 2\sigma_{12} + \sigma_{22} & \sigma_{11} - \sigma_{22} \\ \sigma_{11} - \sigma_{22} & \sigma_{11} + 2\sigma_{12} + \sigma_{22} \end{bmatrix} \end{aligned}$$

Note that if  $\sigma_{11} = \sigma_{22}$ , that is, if  $X_1$  and  $X_2$  have equal variances, the off-diagonal terms in  $\boldsymbol{\Sigma}_Z$  vanish. This demonstrates the well-known result that the sum and difference of two random variables with identical variances are uncorrelated. ■

#### Partitioning the Sample Mean Vector and Covariance Matrix

Many of the matrix results in this section have been expressed in terms of population means and variances (covariances). Results in Equations (2.36),



(2-37), (2-38), and (2-40) also hold if the population quantities are replaced by their appropriately defined sample counterparts.

Let  $\bar{x}' = [\bar{x}_1, \bar{x}_2, \dots, \bar{x}_p]$  be the vector of sample averages constructed from  $n$  observations on  $p$  variables  $X_1, X_2, \dots, X_p$ , and let

$$S_n = \begin{bmatrix} s_{11} & \dots & s_{1p} \\ \vdots & \ddots & \vdots \\ s_{p1} & \dots & s_{pp} \end{bmatrix} = \begin{bmatrix} \frac{1}{n} \sum_{j=1}^n (x_{1j} - \bar{x}_1)^2 & \dots & \frac{1}{n} \sum_{j=1}^n (x_{1j} - \bar{x}_1)(x_{pj} - \bar{x}_p) \\ \vdots & \ddots & \vdots \\ \frac{1}{n} \sum_{j=1}^n (x_{pj} - \bar{x}_p)(x_{1j} - \bar{x}_1) & \dots & \frac{1}{n} \sum_{j=1}^n (x_{pj} - \bar{x}_p)^2 \end{bmatrix} \quad (2-46)$$

be the corresponding sample variance-covariance matrix.

The sample mean vector and covariance matrix can be partitioned in order to distinguish quantities corresponding to groups of variables. Thus

$$\bar{x}_{(p \times 1)} = \begin{bmatrix} \bar{x}_1 \\ \vdots \\ \bar{x}_q \\ \bar{x}_{q+1} \\ \vdots \\ \bar{x}_p \end{bmatrix} = \begin{bmatrix} \bar{x}^{(1)} \\ \bar{x}^{(2)} \end{bmatrix} \quad (2-47)$$

and

$$S_n = \begin{bmatrix} s_{11} & \dots & s_{1q} & \dots & s_{1,q+1} & \dots & s_{1p} \\ \vdots & \ddots & \vdots & \ddots & \vdots & \ddots & \vdots \\ s_{q1} & \dots & s_{qq} & \dots & s_{q,q+1} & \dots & s_{qp} \\ \vdots & \ddots & \vdots & \ddots & \vdots & \ddots & \vdots \\ s_{q+1,1} & \dots & s_{q+1,q} & \dots & s_{q+1,q+1} & \dots & s_{q+1,p} \\ \vdots & \ddots & \vdots & \ddots & \vdots & \ddots & \vdots \\ s_{p1} & \dots & s_{pq} & \dots & s_{p,q+1} & \dots & s_{pp} \end{bmatrix}$$

$$= \begin{bmatrix} s_{11} & \dots & s_{1q} \\ \vdots & \ddots & \vdots \\ s_{q1} & \dots & s_{qq} \\ \vdots & \ddots & \vdots \\ s_{q+1,1} & \dots & s_{q+1,q} \\ \vdots & \ddots & \vdots \\ s_{p1} & \dots & s_{pq} \end{bmatrix} \quad (2-48)$$

where  $\bar{x}^{(1)}$  and  $\bar{x}^{(2)}$  are the sample mean vectors constructed from observations  $x^{(1)} = [x_1, \dots, x_q]'$  and  $x^{(2)} = [x_{q+1}, \dots, x_p]'$ , respectively;  $S_{11}$  is the sample

covariance matrix computed from observations  $x^{(1)}$ ;  $S_{22}$  is the sample covariance matrix computed from observations  $x^{(2)}$ ; and  $S_{21} = S_{12}$  is the sample covariance matrix for elements of  $x^{(1)}$  and elements of  $x^{(2)}$ .

## 2.7 MATRIX INEQUALITIES AND MAXIMIZATION

Maximization principles play an important role in several multivariate techniques. Linear discriminant analysis, for example, is concerned with allocating observations to predetermined groups. The allocation rule is often a linear function of measurements that maximizes the separation between groups relative to their within-group variability. As another example, principal components are linear combinations of measurements with maximum variability.

The matrix inequalities presented in this section will easily allow us to derive certain maximization results, which will be referenced in later chapters.

**Cauchy-Schwarz Inequality.** Let  $b$  and  $d$  be any two  $(p \times 1)$  vectors. Then

$$(b'd)^2 \leq (b'b)(d'd) \quad (2-48)$$

with equality if and only if  $b = cd$  (or  $d = cb$ ) from some constant  $c$ .

*Proof.* The inequality is obvious if either  $b = 0$  or  $d = 0$ . Excluding this possibility, consider the vector  $b - xd$ , where  $x$  is an arbitrary scalar. Since the length of  $b - xd$  is positive for  $b - xd \neq 0$ , in this case

$$0 < (b - xd)'(b - xd) = b'b - xd'b - b'd'xd + x^2d'd \\ = b'b - 2x(b'd) + x^2(d'd)$$

The expression immediately above is quadratic in  $x$ . If we complete the square by adding and subtracting the scalar  $(b'd)^2/d'd$ , we get

$$0 < b'b - \frac{(b'd)^2}{d'd} + \frac{(b'd)^2}{d'd} - 2x(b'd) + x^2(d'd) \\ = b'b - \frac{(b'd)^2}{d'd} + (d'd) \left[ x - \frac{b'd}{d'd} \right]^2$$

The term in brackets is zero if we choose  $x = b'd/d'd$ , so we conclude

$$0 < b'b - \frac{(b'd)^2}{d'd}$$

or  $(b'd)^2 < (b'b)(d'd)$  if  $b \neq xd$  for some  $x$ .

Note that if  $b = cd$ ,  $0 = (b - cd)'(b - cd)$  and the same argument produces  $(b'd)^2 = (b'b)(d'd)$ . ■

A simple but important extension of the Cauchy-Schwarz inequality follows directly.

**Extended Cauchy-Schwarz Inequality.** Let  $b$  and  $d$  be any two vectors  $(p \times 1)$  and let  $B$  be a positive definite matrix. Then



with equality if and only if  $\mathbf{b} = c\mathbf{B}^{-1/2}\mathbf{d}$  (or  $\mathbf{d} = c\mathbf{B}\mathbf{b}$ ) for some constant  $c$ . (2-49)

*Proof.* The inequality is obvious when  $\mathbf{b} = \mathbf{0}$  or  $\mathbf{d} = \mathbf{0}$ . For cases other than these, consider the square-root matrix  $\mathbf{B}^{1/2}$  defined in terms of its eigenvalues  $\lambda_i$ , and the normalized eigenvectors,  $\mathbf{e}_i$ , as:  $\mathbf{B}^{1/2} = \sum_{i=1}^p \sqrt{\lambda_i} \mathbf{e}_i \mathbf{e}_i'$ . Setting [see also (2-22)]

$$\mathbf{B}^{-1/2} = \sum_{i=1}^p \frac{1}{\sqrt{\lambda_i}} \mathbf{e}_i \mathbf{e}_i'$$

it follows that

$$\mathbf{b}'\mathbf{d} = \mathbf{b}'\mathbf{B}^{1/2}\mathbf{B}^{-1/2}\mathbf{d} = (\mathbf{B}^{1/2}\mathbf{b})'(\mathbf{B}^{-1/2}\mathbf{d})$$

and the proof is completed by applying the Cauchy-Schwarz inequality to the vectors  $(\mathbf{B}^{1/2}\mathbf{b})$  and  $(\mathbf{B}^{-1/2}\mathbf{d})$ .

The extended Cauchy-Schwarz inequality gives rise to the following maximization result.

**Maximization Lemma.** Let  $\mathbf{B}$  be positive definite and  $\mathbf{d}$  be a given vector. Then for an arbitrary nonzero vector  $\mathbf{x}$ , (2-50)

$$\max_{\mathbf{x} \neq \mathbf{0}} \frac{(\mathbf{x}'\mathbf{d})^2}{\mathbf{x}'\mathbf{B}\mathbf{x}} = \mathbf{d}'\mathbf{B}^{-1}\mathbf{d} \quad (2-50)$$

with the maximum attained when  $\mathbf{x} = c\mathbf{B}^{-1}\mathbf{d}$  for any constant  $c \neq 0$ .

*Proof.* By the extended Cauchy-Schwarz inequality,  $(\mathbf{x}'\mathbf{d})^2 \leq (\mathbf{x}'\mathbf{B}\mathbf{x})(\mathbf{d}'\mathbf{B}^{-1}\mathbf{d})$ . Because  $\mathbf{x} \neq \mathbf{0}$  and  $\mathbf{B}$  is positive definite,  $\mathbf{x}'\mathbf{B}\mathbf{x} > 0$ . Dividing both sides of the inequality by the positive scalar  $\mathbf{x}'\mathbf{B}\mathbf{x}$  yields the upper bound

$$\frac{(\mathbf{x}'\mathbf{d})^2}{\mathbf{x}'\mathbf{B}\mathbf{x}} \leq \mathbf{d}'\mathbf{B}^{-1}\mathbf{d}$$

Taking the maximum over  $\mathbf{x}$  gives Equation (2-50) because the bound is attained for  $\mathbf{x} = c\mathbf{B}^{-1}\mathbf{d}$ .

A final maximization result will provide us with an interpretation of eigenvalues.

**Maximization of Quadratic Forms for Points on the Unit Sphere.** Let  $\mathbf{B}$  be a positive definite matrix with eigenvalues  $\lambda_1 \geq \lambda_2 \geq \dots \geq \lambda_p > 0$  and associated normalized eigenvectors  $\mathbf{e}_1, \mathbf{e}_2, \dots, \mathbf{e}_p$ . Then

$$\begin{aligned} \max_{\mathbf{x} \neq \mathbf{0}} \frac{\mathbf{x}'\mathbf{B}\mathbf{x}}{\mathbf{x}'\mathbf{x}} &= \lambda_1 & \text{attained when } \mathbf{x} = \mathbf{e}_1 \\ \min_{\mathbf{x} \neq \mathbf{0}} \frac{\mathbf{x}'\mathbf{B}\mathbf{x}}{\mathbf{x}'\mathbf{x}} &= \lambda_p & \text{attained when } \mathbf{x} = \mathbf{e}_p \end{aligned} \quad (2-51)$$

Moreover,

$$\max_{\mathbf{x} \in \mathbf{e}_1, \dots, \mathbf{e}_p} \frac{\mathbf{x}'\mathbf{B}\mathbf{x}}{\mathbf{x}'\mathbf{x}} = \lambda_{k+1} \quad \text{attained when } \mathbf{x} = \mathbf{e}_{k+1}, k = 1, 2, \dots, p-1 \quad (2-52)$$

where the symbol  $\perp$  is read "perpendicular to."

*Proof.* Let  $\mathbf{P}$  be the orthogonal matrix whose columns are the eigenvectors  $\mathbf{e}_1, \mathbf{e}_2, \dots, \mathbf{e}_p$  and  $\mathbf{A}$  be the diagonal matrix with eigenvalues  $\lambda_1, \lambda_2, \dots, \lambda_p$  along the main diagonal. Let  $\mathbf{B}^{1/2} = \mathbf{P}\mathbf{A}^{1/2}\mathbf{P}'$  [see (2-22)] and  $\mathbf{y} = \mathbf{P}'\mathbf{x}$ . Consequently,  $\mathbf{x} \neq \mathbf{0}$  implies  $\mathbf{y} \neq \mathbf{0}$ . Thus

$$\begin{aligned} \frac{\mathbf{x}'\mathbf{B}\mathbf{x}}{\mathbf{x}'\mathbf{x}} &= \frac{\mathbf{x}'\mathbf{B}^{1/2}\mathbf{B}^{1/2}\mathbf{x}}{\mathbf{x}'\mathbf{P}\mathbf{P}'\mathbf{x}} = \frac{\mathbf{x}'\mathbf{P}\mathbf{A}^{1/2}\mathbf{P}'\mathbf{P}\mathbf{A}^{1/2}\mathbf{P}'\mathbf{x}}{\mathbf{y}'\mathbf{y}} = \frac{\mathbf{y}'\mathbf{A}\mathbf{y}}{\mathbf{y}'\mathbf{y}} \leq \lambda_1 \frac{\sum_{i=1}^p y_i^2}{\sum_{i=1}^p y_i^2} = \lambda_1 \end{aligned} \quad (2-53)$$

Setting  $\mathbf{x} = \mathbf{e}_1$  gives

$$\mathbf{y} = \mathbf{P}'\mathbf{e}_1 = \begin{bmatrix} 1 \\ 0 \\ \vdots \\ 0 \end{bmatrix}$$

since

$$\mathbf{e}_i'\mathbf{e}_k = \begin{cases} 1 & k = i \\ 0 & k \neq i \end{cases}$$

For this choice of  $\mathbf{x}$ ,  $\mathbf{y}'\mathbf{A}\mathbf{y}/\mathbf{y}'\mathbf{y} = \lambda_1/1 = \lambda_1$ , or

$$\frac{\mathbf{e}_1'\mathbf{B}\mathbf{e}_1}{\mathbf{e}_1'\mathbf{e}_1} = \mathbf{e}_1'\mathbf{B}\mathbf{e}_1 = \lambda_1 \quad (2-54)$$

A similar argument produces the second part of (2-51).

Now  $\mathbf{x} = \mathbf{P}\mathbf{y} = y_1\mathbf{e}_1 + y_2\mathbf{e}_2 + \dots + y_p\mathbf{e}_p$ , so  $\mathbf{x} \perp \mathbf{e}_1, \dots, \mathbf{e}_k$  implies  $0 = \mathbf{e}_i'\mathbf{x} = y_1\mathbf{e}_i'\mathbf{e}_1 + y_2\mathbf{e}_i'\mathbf{e}_2 + \dots + y_p\mathbf{e}_i'\mathbf{e}_p = y_i$ ,  $i \leq k$

Therefore, for  $\mathbf{x}$  perpendicular to the first  $k$  eigenvectors  $\mathbf{e}_i$ , the left-hand side of the inequality in (2-53) becomes

$$\frac{\mathbf{x}'\mathbf{B}\mathbf{x}}{\mathbf{x}'\mathbf{x}} = \frac{\sum_{i=k+1}^p \lambda_i y_i^2}{\sum_{i=k+1}^p y_i^2}$$

Taking  $y_{k+1} = 1, y_{k+2} = \dots = y_p = 0$  gives the asserted maximum. ■

For a fixed  $\mathbf{x}_0 \neq \mathbf{0}$ ,  $\mathbf{x}_0'\mathbf{B}\mathbf{x}_0/\mathbf{x}_0'\mathbf{x}_0$  has the same value as  $\mathbf{x}'\mathbf{B}\mathbf{x}$ , where  $\mathbf{x}' = \mathbf{x}_0/\sqrt{\mathbf{x}_0'\mathbf{x}_0}$  is of unit length. Consequently, Equation (2-51) says that the largest

eigenvalue,  $\lambda_1$ , is the maximum value of the quadratic form  $\mathbf{x}'\mathbf{B}\mathbf{x}$  for all points  $\mathbf{x}$  whose distance from the origin is unity. Similarly,  $\lambda_p$  is the smallest value of the quadratic form for all points  $\mathbf{x}$  one unit from the origin. The largest and smallest eigenvalues thus represent extremal values of  $\mathbf{x}'\mathbf{B}\mathbf{x}$  for points on the unit sphere. The "intermediate" eigenvalues of the  $(p \times p)$  positive definite matrix  $\mathbf{B}$  also have an interpretation as extremal values when  $\mathbf{x}$  is further restricted to be perpendicular to the earlier choices.

## Supplement 2A: Vectors and Matrices: Basic Concepts

### Vectors

Many concepts, such as a person's health, intellectual abilities, or personality, cannot be adequately quantified as a single number. Rather, several different measurements  $x_1, x_2, \dots, x_m$  are required.

**Definition 2A.1.** An  $m$ -tuple of real numbers  $(x_1, x_2, \dots, x_i, \dots, x_m)$  arranged in a column is called a *vector* and is denoted by a boldfaced, lowercase letter.

Examples of vectors are

$$\mathbf{x} = \begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_m \end{bmatrix}, \quad \mathbf{a} = \begin{bmatrix} 1 \\ 0 \\ 0 \end{bmatrix}, \quad \mathbf{b} = \begin{bmatrix} 1 \\ -1 \\ 1 \end{bmatrix}, \quad \mathbf{y} = \begin{bmatrix} 1 \\ 2 \\ -2 \end{bmatrix}$$

Vectors are said to be equal if their corresponding entries are the same.

**Definition 2A.2 (Scalar Multiplication).** Let  $c$  be an arbitrary scalar. Then the product  $c\mathbf{x}$  is a vector with  $i$ th entry  $cx_i$ .

To illustrate scalar multiplication, take  $c_1 = 5$  and  $c_2 = -1.2$ . Then

$$c_1\mathbf{y} = 5 \begin{bmatrix} 1 \\ 2 \\ -2 \end{bmatrix} = \begin{bmatrix} 5 \\ 10 \\ -10 \end{bmatrix} \quad \text{and} \quad c_2\mathbf{y} = (-1.2) \begin{bmatrix} 1 \\ 2 \\ -2 \end{bmatrix} = \begin{bmatrix} -1.2 \\ -2.4 \\ 2.4 \end{bmatrix}$$

**Definition 2A.3 (Vector Addition).** The sum of two vectors  $\mathbf{x}$  and  $\mathbf{y}$ , each having the same number of entries, is that vector

$$\mathbf{z} = \mathbf{x} + \mathbf{y} \quad \text{with } i\text{th entry } z_i = x_i + y_i.$$

Thus

$$\begin{bmatrix} 3 \\ -1 \\ 4 \end{bmatrix} + \begin{bmatrix} 1 \\ 2 \\ -2 \end{bmatrix} = \begin{bmatrix} 4 \\ 1 \\ 2 \end{bmatrix} \quad \mathbf{x} + \mathbf{y} = \mathbf{z}$$

Taking the zero vector,  $\mathbf{0}$ , to be the  $m$ -tuple  $(0, 0, \dots, 0)$  and the vector  $-\mathbf{x}$  to be the  $m$ -tuple  $(-x_1, -x_2, \dots, -x_m)$ , the two operations of scalar multiplication and vector addition can be combined in a useful manner.

**Definition 2A.4.** The space of all real  $m$ -tuples, with scalar multiplication and vector addition as defined above, is called a *vector space*.



**Definition 2A.5.** The vector  $y = a_1x_1 + a_2x_2 + \dots + a_nx_n$  is a *linear combination* of the vectors  $x_1, x_2, \dots, x_n$ . The set of all linear combinations of  $x_1, x_2, \dots, x_n$  is called their *linear span*.

**Definition 2A.6.** A set of vectors  $x_1, x_2, \dots, x_n$  is said to be *linearly dependent* if there exist  $k$  numbers  $(a_1, a_2, \dots, a_n)$ , not all zero, such that

$$a_1x_1 + a_2x_2 + \dots + a_nx_n = 0$$

Otherwise the set of vectors is said to be *linearly independent*.

If one of the vectors, for example,  $x_1$ , is 0, the set is linearly dependent (let  $a_1$  be the only nonzero coefficient in Definition 2A.6).

The familiar vectors with a one as an entry and zeros elsewhere are linearly independent. For  $m = 4$

$$x_1 = \begin{bmatrix} 1 \\ 0 \\ 0 \\ 0 \end{bmatrix}, \quad x_2 = \begin{bmatrix} 0 \\ 1 \\ 0 \\ 0 \end{bmatrix}, \quad x_3 = \begin{bmatrix} 0 \\ 0 \\ 1 \\ 0 \end{bmatrix}, \quad x_4 = \begin{bmatrix} 0 \\ 0 \\ 0 \\ 1 \end{bmatrix}$$

so

$$0 = a_1x_1 + a_2x_2 + a_3x_3 + a_4x_4 = \begin{bmatrix} a_1 \cdot 1 + a_2 \cdot 0 + a_3 \cdot 0 + a_4 \cdot 0 \\ a_1 \cdot 0 + a_2 \cdot 1 + a_3 \cdot 0 + a_4 \cdot 0 \\ a_1 \cdot 0 + a_2 \cdot 0 + a_3 \cdot 1 + a_4 \cdot 0 \\ a_1 \cdot 0 + a_2 \cdot 0 + a_3 \cdot 0 + a_4 \cdot 1 \end{bmatrix} = \begin{bmatrix} a_1 \\ a_2 \\ a_3 \\ a_4 \end{bmatrix}$$

implies  $a_1 = a_2 = a_3 = a_4 = 0$ .

As another example, let  $k = 3$  and  $m = 3$ .

$$x_1 = \begin{bmatrix} 1 \\ 1 \\ 1 \end{bmatrix}, \quad x_2 = \begin{bmatrix} 2 \\ 5 \\ -1 \end{bmatrix}, \quad x_3 = \begin{bmatrix} 0 \\ 1 \\ -1 \end{bmatrix}$$

Then

$$2x_1 - x_2 + 3x_3 = 0$$

Thus  $x_1, x_2, x_3$  are a linearly dependent set of vectors since any one can be written as a linear combination of the others (for example,  $x_2 = 2x_1 + 3x_3$ ).

**Definition 2A.7.** Any set of  $m$  linearly independent vectors is called a *basis* for the vector space of all  $m$ -tuples of real numbers.

**Result 2A.1.** Every vector can be expressed as a unique linear combination of a fixed basis.

With  $m = 4$ , the usual choice of a basis is

$$\begin{bmatrix} 1 \\ 0 \\ 0 \\ 0 \end{bmatrix}, \begin{bmatrix} 0 \\ 1 \\ 0 \\ 0 \end{bmatrix}, \begin{bmatrix} 0 \\ 0 \\ 1 \\ 0 \end{bmatrix}, \begin{bmatrix} 0 \\ 0 \\ 0 \\ 1 \end{bmatrix}$$

where these four vectors were shown to be linearly independent above. Any vector  $x$  can be uniquely expressed as

$$\begin{bmatrix} x_1 \\ x_2 \\ x_3 \\ x_4 \end{bmatrix} = \begin{bmatrix} 1 \\ 0 \\ 0 \\ 0 \end{bmatrix} + x_2 \begin{bmatrix} 0 \\ 1 \\ 0 \\ 0 \end{bmatrix} + x_3 \begin{bmatrix} 0 \\ 0 \\ 1 \\ 0 \end{bmatrix} + x_4 \begin{bmatrix} 0 \\ 0 \\ 0 \\ 1 \end{bmatrix} = \begin{bmatrix} x_1 \\ x_2 \\ x_3 \\ x_4 \end{bmatrix} = x$$

A vector consisting of  $m$  elements may be regarded geometrically as a point in  $m$ -dimensional space. For example, with  $m = 2$ , the vector  $x$  may be regarded as representing the point in the plane with coordinates  $x_1$  and  $x_2$ .



Vectors have the geometrical properties of length and direction.

**Definition 2A.8.** The *length* of a vector of  $m$  elements emanating from the origin is given by the Pythagorean formula:

$$\text{length of } x = L_x = \sqrt{x_1^2 + x_2^2 + \dots + x_m^2}$$

**Definition 2A.9.** The *angle*  $\theta$  between two vectors  $x$  and  $y$ , both having  $m$  entries, is defined as follows.

$$\cos(\theta) = \frac{(x_1y_1 + x_2y_2 + \dots + x_my_m)}{L_xL_y}$$

where  $L_x$  = length of  $x$  and  $L_y$  = length of  $y$ ,  $x_1, x_2, \dots, x_m$  are the elements of  $x$ , and  $y_1, y_2, \dots, y_m$  are the elements of  $y$ .

Let

$$x = \begin{bmatrix} -1 \\ 5 \\ 2 \\ -2 \end{bmatrix} \quad \text{and} \quad y = \begin{bmatrix} 4 \\ -3 \\ 0 \\ 1 \end{bmatrix}$$

Then the length of  $x$ , length of  $y$ , and the angle between the two vectors are

$$\text{length of } x = \sqrt{(-1)^2 + 5^2 + 2^2 + (-2)^2} = \sqrt{34} = 5.83$$

$$\text{length of } y = \sqrt{4^2 + (-3)^2 + 0^2 + 1^2} = \sqrt{26} = 5.10$$

and

$$\cos(\theta) = \frac{1}{L_xL_y} [x_1y_1 + x_2y_2 + x_3y_3 + x_4y_4]$$

$$\begin{aligned}
 &= \frac{1}{\sqrt{34}\sqrt{26}}[(-1)(4) + 5(-3) + 2(0) + (-2)(1)] \\
 &= \frac{1}{5.83 \times 5.10}[-21] = -.706
 \end{aligned}$$

Consequently,  $\theta = 135^\circ$ .

**Definition 2A.10.** The *inner (or dot) product* of two vectors  $\mathbf{x}$  and  $\mathbf{y}$  with the same number of entries is defined as the sum of component products

$$x_1y_1 + x_2y_2 + \cdots + x_my_m$$

We use the notation  $\mathbf{x}'\mathbf{y}$  or  $\mathbf{y}'\mathbf{x}$  to denote this inner product.

With the  $\mathbf{x}'\mathbf{y}$  notation, we may express the length of a vector and the cosine of the angle between two vectors as

$$\begin{aligned}
 L_x &= \text{length of } \mathbf{x} = \sqrt{x_1^2 + x_2^2 + \cdots + x_m^2} = \sqrt{\mathbf{x}'\mathbf{x}} \\
 \cos(\theta) &= \frac{\mathbf{x}'\mathbf{y}}{\sqrt{\mathbf{x}'\mathbf{x}}\sqrt{\mathbf{y}'\mathbf{y}}}
 \end{aligned}$$

**Definition 2A.11.** When the angle between two vectors  $\mathbf{x}$ ,  $\mathbf{y}$  is  $\theta = 90^\circ$  or  $270^\circ$ , we say that  $\mathbf{x}$  and  $\mathbf{y}$  are perpendicular. Since  $\cos(\theta) = 0$  only if  $\theta = 90^\circ$  or  $270^\circ$ , the condition becomes

We write  $\mathbf{x} \perp \mathbf{y}$ .

The basis vectors

$$\begin{bmatrix} 1 \\ 0 \\ 0 \\ 0 \end{bmatrix}, \begin{bmatrix} 0 \\ 1 \\ 0 \\ 0 \end{bmatrix}, \begin{bmatrix} 0 \\ 0 \\ 1 \\ 0 \end{bmatrix}, \begin{bmatrix} 0 \\ 0 \\ 0 \\ 1 \end{bmatrix}$$

are mutually perpendicular. Also, each has length one. The same construction holds for any number of entries  $m$ .

**Result 2A.2**

- (a)  $\mathbf{z}$  is perpendicular to every vector if and only if  $\mathbf{z} = \mathbf{0}$ .
- (b) If  $\mathbf{z}$  is perpendicular to each vector  $\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_k$ , then  $\mathbf{z}$  is perpendicular to their linear span.
- (c) Mutually perpendicular vectors are linearly independent.

**Definition 2A.12.** The *projection (or shadow)* of a vector  $\mathbf{x}$  on a vector  $\mathbf{y}$  is

$$\text{projection of } \mathbf{x} \text{ on } \mathbf{y} = \frac{(\mathbf{x}'\mathbf{y})}{L_y^2} \mathbf{y}$$

If  $\mathbf{y}$  has unit length so that  $L_y = 1$ ,

$$\text{projection of } \mathbf{x} \text{ on } \mathbf{y} = (\mathbf{x}'\mathbf{y})\mathbf{y}$$

**Result 2A.3 (Gram-Schmidt Process).** Given linearly independent vectors  $\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_k$ , there exist mutually perpendicular vectors  $\mathbf{u}_1, \mathbf{u}_2, \dots, \mathbf{u}_k$  with the same linear span. These may be constructed sequentially by setting

$$\begin{aligned}
 \mathbf{u}_1 &= \mathbf{x}_1 \\
 \mathbf{u}_2 &= \mathbf{x}_2 - \frac{(\mathbf{x}_2'\mathbf{u}_1)}{\mathbf{u}_1'\mathbf{u}_1} \mathbf{u}_1 \\
 &\vdots \\
 \mathbf{u}_k &= \mathbf{x}_k - \frac{(\mathbf{x}_k'\mathbf{u}_1)}{\mathbf{u}_1'\mathbf{u}_1} \mathbf{u}_1 - \cdots - \frac{(\mathbf{x}_k'\mathbf{u}_{k-1})}{\mathbf{u}_{k-1}'\mathbf{u}_{k-1}} \mathbf{u}_{k-1}
 \end{aligned}$$

We can also convert the  $\mathbf{u}$ 's to unit length by setting  $\mathbf{z}_j = \mathbf{u}_j/\sqrt{\mathbf{u}_j'\mathbf{u}_j}$ . In this construction,  $(\mathbf{x}_i'\mathbf{z}_j)\mathbf{z}_j$  is the projection of  $\mathbf{x}_i$  on  $\mathbf{z}_j$  and  $\sum_{j=1}^k (\mathbf{x}_i'\mathbf{z}_j)\mathbf{z}_j$  is the projection of  $\mathbf{x}_i$  on the linear span of  $\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_{k-1}$ . ■

For example, to construct perpendicular vectors from

$$\mathbf{x}_1 = \begin{bmatrix} 4 \\ 0 \\ 0 \\ 2 \end{bmatrix} \quad \text{and} \quad \mathbf{x}_2 = \begin{bmatrix} 3 \\ 1 \\ 0 \\ -1 \end{bmatrix}$$

we take

$$\mathbf{u}_1 = \mathbf{x}_1 = \begin{bmatrix} 4 \\ 0 \\ 0 \\ 2 \end{bmatrix}$$

so

$$\mathbf{u}_1'\mathbf{u}_1 = 4^2 + 0^2 + 0^2 + 2^2 = 20$$

and

$$\mathbf{x}_2'\mathbf{u}_1 = 3(4) + 1(0) + 0(0) + 1(2) = 10$$

Thus

$$\mathbf{u}_2 = \begin{bmatrix} 3 \\ 1 \\ -10 \\ -2 \end{bmatrix} - \frac{10}{20} \begin{bmatrix} 4 \\ 0 \\ 0 \\ 2 \end{bmatrix} = \begin{bmatrix} 1 \\ 1 \\ 0 \\ -2 \end{bmatrix} \quad \text{and} \quad \mathbf{z}_1 = \frac{1}{\sqrt{20}} \begin{bmatrix} 4 \\ 0 \\ 0 \\ 2 \end{bmatrix}, \quad \mathbf{z}_2 = \frac{1}{\sqrt{6}} \begin{bmatrix} 1 \\ 1 \\ 0 \\ -2 \end{bmatrix}$$

**Matrices**

**Definition 2A.13.** An  $m \times k$  matrix, generally denoted by a boldface uppercase letter such as  $\mathbf{A}$ ,  $\mathbf{R}$ ,  $\mathbf{\Sigma}$ , and so forth, is a rectangular array of elements having  $m$  rows and  $k$  columns.



Examples of matrices are

$$A = \begin{bmatrix} -7 & 2 \\ 0 & 1 \\ 3 & 4 \end{bmatrix}, \quad B = \begin{bmatrix} x & 3 & 0 \\ 4 & -2 & 1/x \end{bmatrix}, \quad I = \begin{bmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \end{bmatrix}$$

$$\Sigma = \begin{bmatrix} 1 & .7 & -.3 \\ .7 & 2 & 1 \\ -.3 & 1 & 8 \end{bmatrix}, \quad E = [e_i]$$

In our work, the matrix elements will be real numbers or functions taking on values in the real numbers.

**Definition 2A.14.** The *dimension* (abbreviated *dim*) of an  $m \times k$  matrix is the ordered pair  $(m, k)$ ;  $m$  is the row dimension and  $k$  is the column dimension. The dimension of a matrix is frequently indicated in parentheses below the letter representing the matrix. Thus the  $m \times k$  matrix  $A$  is denoted by  $A_{(m \times k)}$ . In the examples above the dimension of the matrix  $\Sigma$  is  $3 \times 3$  and this information can be conveyed by writing  $\Sigma_{(3 \times 3)}$ .

An  $m \times k$  matrix, call it  $A$ , of arbitrary constants can be written

$$A_{(m \times k)} = \begin{bmatrix} a_{11} & a_{12} & \cdots & a_{1k} \\ a_{21} & a_{22} & \cdots & a_{2k} \\ \vdots & \vdots & \ddots & \vdots \\ a_{m1} & a_{m2} & \cdots & a_{mk} \end{bmatrix}$$

or more compactly as  $A_{(m \times k)} = \{a_{ij}\}$ , where the index  $i$  refers to the row and the index  $j$  refers to the column.

An  $m \times 1$  matrix is referred to as a *column vector*. A  $1 \times k$  matrix is referred to as a *row vector*. Since matrices can be considered as vectors side by side, it is natural to define multiplication by a scalar and the addition of two matrices with the same dimensions.

**Definition 2A.15.** Two matrices  $A_{(m \times k)} = \{a_{ij}\}$  and  $B_{(m \times k)} = \{b_{ij}\}$  are said to be *equal*, written  $A = B$ , if  $a_{ij} = b_{ij}$ ,  $i = 1, 2, \dots, m$ ,  $j = 1, 2, \dots, k$ . That is, two matrices are equal if:

- Their dimensionality is the same.
- Every corresponding element is the same.

**Definition 2A.16 (Matrix Addition).** Let the matrices  $A$  and  $B$  both be of  $\dim m \times k$  with arbitrary elements  $a_{ij}$  and  $b_{ij}$ ,  $i = 1, 2, \dots, m$ ,  $j = 1, 2, \dots, k$ , respectively. The sum of the matrices  $A$  and  $B$  is an  $m \times k$  matrix  $C$ , written  $C = A + B$ , such that the arbitrary element of  $C$ ,  $c_{ij}$ , is given by

$$c_{ij} = a_{ij} + b_{ij} \quad i = 1, 2, \dots, m, \quad j = 1, 2, \dots, k$$

Note that the addition of matrices is defined only for matrices of the same dimension.

For example

$$\begin{bmatrix} 3 & 2 & 3 \\ 4 & 1 & 1 \end{bmatrix} + \begin{bmatrix} 3 & 6 & 7 \\ 2 & -1 & 0 \end{bmatrix} = \begin{bmatrix} 6 & 8 & 10 \\ 6 & 0 & 1 \end{bmatrix}$$

$$A + B = C$$

**Definition 2A.17 (Scalar Multiplication).** Let  $c$  be an arbitrary scalar and  $A_{(m \times k)} = \{a_{ij}\}$ . Then  $cA_{(m \times k)} = Ac_{(m \times k)} = B_{(m \times k)} = \{b_{ij}\}$ , where  $b_{ij} = ca_{ij} = a_{ij}c$ ,  $i = 1, 2, \dots, m$ ,  $j = 1, 2, \dots, k$ .

Multiplication of a matrix by a scalar produces a new matrix whose elements are the elements of the original matrix, *each* multiplied by the scalar.

For example, if  $c = 2$ ,

$$2 \begin{bmatrix} 3 & -4 \\ 2 & 6 \\ 0 & 5 \end{bmatrix} = \begin{bmatrix} 3 & -4 \\ 2 & 6 \\ 0 & 5 \end{bmatrix} 2 = \begin{bmatrix} 6 & -8 \\ 4 & 12 \\ 0 & 10 \end{bmatrix}$$

$$cA = Ac = B$$

**Definition 2A.18 (Matrix Subtraction).** Let  $A_{(m \times k)} = \{a_{ij}\}$  and  $B_{(m \times k)} = \{b_{ij}\}$  be two matrices of equal dimension. Then the difference between  $A$  and  $B$ , written  $A - B$ , is an  $m \times k$  matrix  $C = \{c_{ij}\}$  given by

$$C = A - B = A + (-1)B$$

That is,  $c_{ij} = a_{ij} + (-1)b_{ij} = a_{ij} - b_{ij}$ ,  $i = 1, 2, \dots, m$ ,  $j = 1, 2, \dots, k$ .

**Definition 2A.19.** Consider the  $m \times k$  matrix  $A$  with arbitrary elements  $a_{ij}$ ,  $i = 1, 2, \dots, m$ ,  $j = 1, 2, \dots, k$ . The *transpose* of the matrix  $A$ , denoted by  $A'$ , is the  $k \times m$  matrix with elements  $a_{ji}$ ,  $j = 1, 2, \dots, k$ ,  $i = 1, 2, \dots, m$ . That is, the transpose of the matrix  $A$  is obtained from  $A$  by interchanging the rows and columns.

As an example, if

$$A_{(2 \times 3)} = \begin{bmatrix} 2 & 1 & 3 \\ 7 & -4 & 6 \end{bmatrix}, \quad \text{then} \quad A'_{(3 \times 2)} = \begin{bmatrix} 2 & 7 \\ 1 & -4 \\ 3 & 6 \end{bmatrix}$$

**Result 2A.4.** For all matrices  $A$ ,  $B$ , and  $C$  (of equal dimension) and scalars  $c$  and  $d$ , the following hold:

- $(A + B) + C = A + (B + C)$
- $A + B = B + A$
- $c(A + B) = cA + cB$
- $(c + d)A = cA + dA$
- $(A + B)' = A' + B'$  (That is, the transpose of the sum is equal to the sum of the transposes.)



- (f)  $(cd)A = c(dA)$   
 (g)  $(cA)' = cA'$

**Definition 2A.20.** If an arbitrary matrix  $A$  has the same number of rows and columns, then  $A$  is called a *square* matrix. In the matrices given above, Definition 2A.13, the matrices  $\Sigma$ ,  $I$ , and  $E$  are square matrices.

**Definition 2A.21.** Let  $A$  be a  $k \times k$  (square) matrix. The matrix  $A$  is said to be *symmetric* if  $A = A'$ . That is,  $A$  is symmetric if  $a_{ij} = a_{ji}$ ,  $i, j = 1, 2, \dots, k$ .

Examples of symmetric matrices are

$$I_{(3 \times 3)} = \begin{bmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \end{bmatrix}, \quad A_{(2 \times 2)} = \begin{bmatrix} 2 & 4 \\ 4 & 1 \end{bmatrix}, \quad B_{(4 \times 4)} = \begin{bmatrix} a & c & e & f \\ c & b & g & d \\ e & g & c & a \\ f & d & a & d \end{bmatrix}$$

**Definition 2A.22.** The  $k \times k$  identity matrix, denoted by  $I_{(k \times k)}$ , is the square matrix with ones on the main (NW-SE) diagonal and zeros elsewhere. The  $3 \times 3$  identity matrix is included in the collection above.

**Definition 2A.23 (Matrix Multiplication).** The product  $AB$  of an  $m \times k$  matrix  $A = \{a_{ij}\}$  and an  $n \times k$  matrix  $B = \{b_{ij}\}$  is the  $m \times k$  matrix  $C$  whose element  $c_{ij}$  is given by

$$c_{ij} = \sum_{k=1}^n a_{ik}b_{kj} \quad i = 1, 2, \dots, m, \quad j = 1, 2, \dots, k$$

Note that for the product  $AB$  to be defined, the column dimension of  $A$  must equal the row dimension of  $B$ . Then the row dimension of  $AB$  equals the row dimension of  $A$  and the column dimension of  $AB$  equals the column dimension of  $B$ .

For example let

$$A_{(2 \times 3)} = \begin{bmatrix} 3 & -1 & 2 \\ 4 & 0 & 5 \end{bmatrix} \quad \text{and} \quad B_{(3 \times 2)} = \begin{bmatrix} 3 & 4 \\ 6 & -2 \\ 4 & 3 \end{bmatrix}$$

Then

$$\begin{bmatrix} 3 & -1 & 2 \\ 4 & 0 & 5 \end{bmatrix}_{(2 \times 3)} \begin{bmatrix} 3 & 4 \\ 6 & -2 \\ 4 & 3 \end{bmatrix}_{(3 \times 2)} = \begin{bmatrix} 11 & 20 \\ 32 & 31 \end{bmatrix}_{(2 \times 2)} = \begin{bmatrix} c_{11} & c_{12} \\ c_{21} & c_{22} \end{bmatrix}$$

where

$$\begin{aligned} c_{11} &= (3)(3) + (-1)(6) + (2)(4) = 11 \\ c_{12} &= (3)(4) + (-1)(-2) + (2)(3) = 20 \\ c_{21} &= (4)(3) + (0)(6) + (5)(4) = 32 \\ c_{22} &= (4)(4) + (0)(-2) + (5)(3) = 31 \end{aligned}$$

As an additional example, consider the product of two vectors. Let

$$x = \begin{bmatrix} 1 \\ 0 \\ -2 \\ 3 \end{bmatrix} \quad \text{and} \quad y = \begin{bmatrix} 2 \\ -3 \\ -1 \\ -8 \end{bmatrix}$$

Then  $x' = [1 \ 0 \ -2 \ 3]$  and

$$x'y = [1 \ 0 \ -2 \ 3] \begin{bmatrix} 2 \\ -3 \\ -1 \\ -8 \end{bmatrix} = [-20] = [2 \ -3 \ -1 \ -8] \begin{bmatrix} 1 \\ 0 \\ -2 \\ 3 \end{bmatrix} = y'x$$

Note that the product  $xy$  is undefined since  $x$  is a  $4 \times 1$  matrix and  $y$  is a  $4 \times 1$  matrix so the column dimension of  $x$ , 1, is unequal to the row dimension of  $y$ , 4. If  $x$  and  $y$  are vectors of the same dimension, such as  $n \times 1$ , both of the products  $xy$  and  $xy'$  are defined. In particular,  $y'x = x'y = x_1y_1 + x_2y_2 + \dots + x_ny_n$  and  $xy'$  is an  $n \times n$  matrix with  $i, j$ th element  $x_iy_j$ .

**Result 2A.5.** For all matrices  $A$ ,  $B$ , and  $C$  (of dimensions such that the indicated products are defined) and a scalar  $c$ ,

- $c(AB) = (cA)B$
- $A(BC) = (AB)C$
- $A(B + C) = AB + AC$
- $(B + C)A = BA + CA$
- $(AB)' = B'A'$

There are several important differences between the algebra of matrices and the algebra of real numbers. Two of these differences follow.

- Matrix multiplication is, in general, not commutative. That is, in general,  $AB \neq BA$ .

Several examples will illustrate the failure of the commutative law (for matrices).

$$\begin{bmatrix} 3 & -1 \\ 4 & 7 \end{bmatrix} \begin{bmatrix} 0 \\ 2 \end{bmatrix} = \begin{bmatrix} -2 \\ 14 \end{bmatrix}$$

but

$$\begin{bmatrix} 0 & 3 \\ 2 & 4 \end{bmatrix} \begin{bmatrix} 3 \\ 7 \end{bmatrix}$$

is not defined.

$$\begin{bmatrix} 1 & 0 & 1 \\ 2 & -3 & 6 \end{bmatrix} \begin{bmatrix} 7 & 6 \\ -3 & 1 \end{bmatrix} = \begin{bmatrix} 9 & 10 \\ 35 & 33 \end{bmatrix}$$

but

$$\begin{bmatrix} 7 & 6 \\ -3 & 1 \end{bmatrix} \begin{bmatrix} 1 & 0 & 1 \\ 2 & -3 & 6 \end{bmatrix} = \begin{bmatrix} 19 & -18 & 43 \\ -1 & -3 & 3 \end{bmatrix}$$



Also

$$\begin{bmatrix} 4 & -1 \\ 0 & 1 \end{bmatrix} \begin{bmatrix} 2 & 1 \\ -3 & 4 \end{bmatrix} = \begin{bmatrix} 11 & 0 \\ -3 & 4 \end{bmatrix}$$

but

$$\begin{bmatrix} 2 & 1 \\ -3 & 4 \end{bmatrix} \begin{bmatrix} 4 & -1 \\ 0 & 1 \end{bmatrix} = \begin{bmatrix} 8 & -1 \\ -12 & 7 \end{bmatrix}$$

2. Let  $\mathbf{0}$  denote the zero matrix, that is, the matrix with zero for every element. In the algebra of real numbers, if the product of two numbers,  $ab$ , is zero, then  $a = 0$  or  $b = 0$ . In matrix algebra, however, the product of two nonzero matrices may be the zero matrix. Hence,

$$\begin{matrix} \mathbf{AB} = \mathbf{0} \\ (m \times n)(n \times k) & (m \times k) \end{matrix}$$

does not imply that  $\mathbf{A} = \mathbf{0}$  or  $\mathbf{B} = \mathbf{0}$ . For example,

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

It is true, however, that if either  $\mathbf{A} = \mathbf{0}$  or  $\mathbf{B} = \mathbf{0}$ , then

$$\begin{matrix} \mathbf{AB} = \mathbf{0} \\ (m \times n)(n \times k) & (m \times k) \end{matrix}$$

**Definition 2A.24.** The *determinant* of the square  $k \times k$  matrix  $\mathbf{A} = [a_{ij}]$  denoted by  $|\mathbf{A}|$ , is the scalar

$$|\mathbf{A}| = a_{11} \quad \text{if } k = 1$$

$$|\mathbf{A}| = \sum_{j=1}^k a_{1j} |\mathbf{A}_{1j}| (-1)^{1+j} \quad \text{if } k > 1$$

where  $\mathbf{A}_{1j}$  is the  $(k-1) \times (k-1)$  matrix obtained by deleting the first row and  $j$ th column of  $\mathbf{A}$ . Also,  $|\mathbf{A}| = \sum_{j=1}^k a_{ij} |\mathbf{A}_{ij}| (-1)^{i+j}$  using the  $i$ th row in place of the first row.

Examples of determinants (evaluated using Definition 2A.24) are

$$\begin{vmatrix} 1 & 3 \\ 6 & 4 \end{vmatrix} = 1(4)(-1)^2 + 3(6)(-1)^3 = 1(4) + 3(6)(-1) = -14$$

In general,

$$\begin{vmatrix} a_{11} & a_{12} \\ a_{21} & a_{22} \end{vmatrix} = a_{11}a_{22}(-1)^2 + a_{12}a_{21}(-1)^3 = a_{11}a_{22} - a_{12}a_{21}$$

$$\begin{vmatrix} 3 & 1 & 6 \\ 7 & 4 & 5 \\ 2 & -7 & 1 \end{vmatrix} = 3 \begin{vmatrix} 4 & 5 \\ -7 & 1 \end{vmatrix} (-1)^2 + 1 \begin{vmatrix} 7 & 5 \\ 2 & 1 \end{vmatrix} (-1)^3 + 6 \begin{vmatrix} 7 & 4 \\ 2 & -7 \end{vmatrix} (-1)^4$$

$$= 3(39) - 1(-3) + 6(-57) = -222$$

$$\begin{vmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \end{vmatrix} = 1 \begin{vmatrix} 1 & 0 \\ 0 & 1 \end{vmatrix} (-1)^2 + 0 \begin{vmatrix} 0 & 0 \\ 0 & 1 \end{vmatrix} (-1)^3 + 0 \begin{vmatrix} 0 & 1 \\ 0 & 0 \end{vmatrix} (-1)^4$$

$$= 1(1) = 1$$

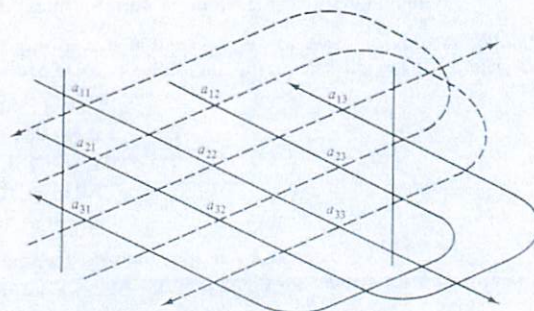
If  $\mathbf{I}$  is the  $k \times k$  identity matrix,  $|\mathbf{I}| = 1$ .

$$\begin{vmatrix} a_{11} & a_{12} & a_{13} \\ a_{21} & a_{22} & a_{23} \\ a_{31} & a_{32} & a_{33} \end{vmatrix} = a_{11} \begin{vmatrix} a_{22} & a_{23} \\ a_{32} & a_{33} \end{vmatrix} (-1)^2 + a_{12} \begin{vmatrix} a_{21} & a_{23} \\ a_{31} & a_{33} \end{vmatrix} (-1)^3$$

$$+ a_{13} \begin{vmatrix} a_{21} & a_{22} \\ a_{31} & a_{32} \end{vmatrix} (-1)^4 = a_{11}a_{22}a_{33} + a_{12}a_{23}a_{31} + a_{13}a_{21}a_{32}$$

$$- a_{31}a_{22}a_{13} - a_{21}a_{12}a_{33} - a_{32}a_{23}a_{11}$$

The determinant of any  $3 \times 3$  matrix can be computed by summing the products of elements along the solid lines and subtracting the products along the dashed lines in the diagram below. This procedure is *not* valid for matrices of higher dimension but, in general, Definition 2A.24 can be employed to evaluate these determinants.



We want to state a result that describes some properties of the determinant. However, we must first introduce some notions related to matrix inverses.

**Definition 2A.25.** The *row rank* of a matrix is the maximum number of linearly independent rows considered as vectors (that is, row vectors). The *column rank* of a matrix is the rank of its set of columns considered as vectors.

For example, let the matrix  $\mathbf{A}$  be

$$\mathbf{A} = \begin{bmatrix} 1 & 1 & 1 \\ 2 & 5 & -1 \\ 0 & 1 & -1 \end{bmatrix}$$

The rows of  $\mathbf{A}$  written as vectors were shown to be linearly dependent after Definition 2A.6. Note that the column rank of  $\mathbf{A}$  is also 2, since



$$\begin{bmatrix} 1 \\ -2 \\ 0 \end{bmatrix} + \frac{1}{5} \begin{bmatrix} 1 \\ -1 \\ 1 \end{bmatrix} + \begin{bmatrix} 1 \\ -1 \\ -1 \end{bmatrix} = \begin{bmatrix} 0 \\ 0 \\ 0 \end{bmatrix}$$

but columns 1 and 2 are linearly independent. This is no coincidence, as the following result indicates.

**Result 2A.6.** The row rank and the column rank of a matrix are equal.

Thus the rank of a matrix is either the row rank or the column rank.

**Definition 2A.26.** A square matrix  $A$  is nonsingular if  $A^{-1}$  exists. If  $A$  is nonsingular, it is called singular. Equivalently, a square matrix is nonsingular if its rank is equal to the number of rows (or columns).

Note that  $Ax = x_1a_1 + x_2a_2 + \dots + x_ka_k$ , where  $a_i$  is the  $i$ th column of  $A$ , so that the condition of nonsingularity is just the statement that the columns of  $A$  are linearly independent.

**Result 2A.7.** Let  $A$  be a nonsingular square matrix of dimension  $k$ . Then there is a unique  $k \times k$  matrix  $B$  such that

$$AB = BA = I$$

where  $I$  is the  $k \times k$  identity matrix.

**Definition 2A.27.** The  $B$  such that  $AB = BA = I$  is called the inverse of  $A$  and is denoted by  $A^{-1}$ . In fact, if  $BA = I$  or  $AB = I$ , then  $B = A^{-1}$  and both products must equal  $I$ .

For example:

$$A = \begin{bmatrix} 2 & 3 \\ 1 & 5 \end{bmatrix} \text{ has } A^{-1} = \begin{bmatrix} \frac{5}{13} & -\frac{3}{13} \\ -\frac{1}{13} & \frac{2}{13} \end{bmatrix}$$

since

$$\begin{bmatrix} 2 & 3 \\ 1 & 5 \end{bmatrix} \begin{bmatrix} \frac{5}{13} & -\frac{3}{13} \\ -\frac{1}{13} & \frac{2}{13} \end{bmatrix} = \begin{bmatrix} \frac{5}{13} & -\frac{3}{13} \\ -\frac{1}{13} & \frac{2}{13} \end{bmatrix} \begin{bmatrix} 2 & 3 \\ 1 & 5 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}$$

**Result 2A.8.**

(a) The inverse of the  $2 \times 2$  matrix

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

is given by

$$A^{-1} = \frac{1}{|A|} \begin{bmatrix} a_{22} & -a_{12} \\ -a_{21} & a_{11} \end{bmatrix}$$

(b) The inverse of a  $3 \times 3$  matrix

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

is given by

$$A^{-1} = \frac{1}{|A|} \begin{bmatrix} \begin{vmatrix} a_{22} & a_{23} \\ a_{32} & a_{33} \end{vmatrix} & -\begin{vmatrix} a_{12} & a_{13} \\ a_{32} & a_{33} \end{vmatrix} & \begin{vmatrix} a_{12} & a_{13} \\ a_{22} & a_{23} \end{vmatrix} \\ -\begin{vmatrix} a_{21} & a_{23} \\ a_{31} & a_{33} \end{vmatrix} & \begin{vmatrix} a_{11} & a_{13} \\ a_{31} & a_{33} \end{vmatrix} & -\begin{vmatrix} a_{11} & a_{13} \\ a_{21} & a_{23} \end{vmatrix} \\ \begin{vmatrix} a_{21} & a_{22} \\ a_{31} & a_{32} \end{vmatrix} & -\begin{vmatrix} a_{11} & a_{12} \\ a_{31} & a_{32} \end{vmatrix} & \begin{vmatrix} a_{11} & a_{12} \\ a_{21} & a_{22} \end{vmatrix} \end{bmatrix}$$

In both cases above, it is clear that  $|A| \neq 0$  if the inverse is to exist.

(c) In general,  $A^{-1}$  has  $j$ ,  $i$ th entry  $[A_{ji}/|A|](i - j)^{n-1}$ , where  $A_{ji}$  is the matrix obtained from  $A$  by deleting the  $i$ th row and  $j$ th column.

**Result 2A.9.** For a square matrix  $A$  of dimension  $k \times k$ , the following are equivalent:

- (a)  $A^{-1}$  exists implies  $|A| \neq 0$  ( $A$  is nonsingular).
- (b)  $|A| \neq 0$ .
- (c) There exists a matrix  $A^{-1}$  such that  $AA^{-1} = A^{-1}A = I$ .

**Result 2A.10.** Let  $A$  and  $B$  be square matrices of the same dimension and let the indicated inverses exist. Then the following hold:

- (a)  $(A^{-1})^{-1} = A$ .
- (b)  $(AB)^{-1} = B^{-1}A^{-1}$ .

The determinant has the following properties.

**Result 2A.11.** Let  $A$  and  $B$  be  $k \times k$  square matrices.

- (a)  $|A| = |A^T|$ .
- (b) If each element of a row (column) of  $A$  is zero, then  $|A| = 0$ .
- (c) If any two rows (columns) of  $A$  are identical, then  $|A| = 0$ .
- (d) If  $A$  is nonsingular, then  $|A| = 1/|A^{-1}|$ ; that is,  $|A||A^{-1}| = 1$ .
- (e)  $|AB| = |A||B|$ .
- (f)  $|cA| = c^k|A|$ , where  $c$  is a scalar.

You are referred to [5] for proofs of parts of Results 2A.9 and 2A.11. Some of these proofs are rather complex and beyond the scope of this book.



**Definition 2A.28.** Let  $A = \{a_{ij}\}$  be a  $k \times k$  square matrix. The *trace* of the matrix  $A$ , written  $\text{tr}(A)$ , is the sum of the diagonal elements; that is,

$$\text{tr}(A) = \sum_{i=1}^k a_{ii}.$$

**Result 2A.12.** Let  $A$  and  $B$  be  $k \times k$  matrices and  $c$  be a scalar.

- (a)  $\text{tr}(cA) = c \text{tr}(A)$
- (b)  $\text{tr}(A \pm B) = \text{tr}(A) \pm \text{tr}(B)$
- (c)  $\text{tr}(AB) = \text{tr}(BA)$
- (d)  $\text{tr}(B^{-1}AB) = \text{tr}(A)$
- (e)  $\text{tr}(AA') = \sum_{i=1}^k \sum_{j=1}^k a_{ij}^2$

**Definition 2A.29.** A square matrix  $A$  is said to be *orthogonal* if its rows, considered as vectors, are mutually perpendicular and have unit lengths; that is,  $AA' = I$ .

**Result 2A.13.** A matrix  $A$  is orthogonal if and only if  $A^{-1} = A'$ . For an orthogonal matrix,  $AA' = A'A = I$ , so the columns are also mutually perpendicular and have unit lengths.

An example of an orthogonal matrix is

$$A = \begin{bmatrix} \frac{1}{\sqrt{2}} & \frac{1}{\sqrt{2}} & 0 & 0 \\ \frac{1}{\sqrt{2}} & -\frac{1}{\sqrt{2}} & 0 & 0 \\ 0 & 0 & \frac{1}{\sqrt{2}} & \frac{1}{\sqrt{2}} \\ 0 & 0 & \frac{1}{\sqrt{2}} & -\frac{1}{\sqrt{2}} \end{bmatrix}$$

Clearly  $A = A'$ , so  $AA' = A'A = AA$ . We verify that  $AA = I = AA' = A'A$  or

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

so  $A' = A^{-1}$  and  $A$  must be an orthogonal matrix.

Square matrices are best understood in terms of quantities called eigenvalues and eigenvectors.

**Definition 2A.30.** Let  $A$  be a  $k \times k$  square matrix and  $I$  be the  $k \times k$  identity matrix. Then the scalars  $\lambda_1, \lambda_2, \dots, \lambda_k$  satisfying the polynomial equation  $|A - \lambda I| = 0$  are called the *eigenvalues* (or *characteristic roots*) of a matrix  $A$ . The equation  $|A - \lambda I| = 0$  (as a function of  $\lambda$ ) is called the *characteristic equation*.

For example, let

$$A = \begin{bmatrix} 1 & 0 \\ 1 & 3 \end{bmatrix}$$

Then

$$|A - \lambda I| = \begin{vmatrix} 1-\lambda & 0 \\ 1 & 3-\lambda \end{vmatrix} = \begin{vmatrix} 1-\lambda & 0 \\ 1 & 3-\lambda \end{vmatrix} = (1-\lambda)(3-\lambda) = 0$$

implies there are two roots,  $\lambda_1 = 1$  and  $\lambda_2 = 3$ . The eigenvalues of  $A$  are 3 and 1. Let

$$A = \begin{bmatrix} 13 & -4 & 2 \\ -4 & 13 & -2 \\ 2 & -2 & 10 \end{bmatrix}$$

Then the equation

$$|A - \lambda I| = \begin{vmatrix} 13-\lambda & -4 & 2 \\ -4 & 13-\lambda & -2 \\ 2 & -2 & 10-\lambda \end{vmatrix} = -\lambda^3 + 36\lambda^2 - 405\lambda + 1458 = 0$$

has three roots:  $\lambda_1 = 9$ ,  $\lambda_2 = 9$ , and  $\lambda_3 = 18$ ; that is, 9, 9, and 18 are the eigenvalues of  $A$ .

**Definition 2A.31.** Let  $A$  be a matrix of dim  $k \times k$  and let  $\lambda$  be an eigenvalue of  $A$ . If  $\mathbf{x}$  is a *nonzero vector* ( $\mathbf{x} \neq \mathbf{0}$ ) such that

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

then  $\mathbf{x}$  is said to be an *eigenvector* (*characteristic vector*) of the matrix  $A$  associated with the *eigenvalue*  $\lambda$ .

An equivalent condition for  $\lambda$  to be a solution of the eigenvalue-eigenvector equation is  $|A - \lambda I| = 0$ . This follows since the statement that  $A\mathbf{x} = \lambda\mathbf{x}$  for some  $\lambda$  and  $\mathbf{x} \neq \mathbf{0}$  implies

$$\mathbf{0} = (A - \lambda I)\mathbf{x} = x_1 \text{col}_1(A - \lambda I) + \dots + x_k \text{col}_k(A - \lambda I)$$

That is, the columns of  $A - \lambda I$  are linearly dependent so, by Result 2A.9(b),  $|A - \lambda I| = 0$  as asserted. Following Definition 2A.30, we have shown that the eigenvalues of

$$A = \begin{bmatrix} 1 & 0 \\ 1 & 3 \end{bmatrix}$$

are  $\lambda_1 = 1$  and  $\lambda_2 = 3$ . The eigenvectors associated with these eigenvalues can be determined by solving the following equations:

$$\begin{bmatrix} 1 & 0 \\ 1 & 3 \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \end{bmatrix} = 1 \begin{bmatrix} x_1 \\ x_2 \end{bmatrix}$$

$$A\mathbf{x} = \lambda_1 \mathbf{x}$$

$$\begin{bmatrix} 1 & 0 \\ 1 & 3 \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \end{bmatrix} = 3 \begin{bmatrix} x_1 \\ x_2 \end{bmatrix}$$

$$A\mathbf{x} = \lambda_2 \mathbf{x}$$

From the first expression,

$$\begin{aligned}x_1 &= x_1 \\x_1 + 3x_2 &= x_2\end{aligned}$$

or

$$x_1 = -2x_2$$

There are many solutions for  $x_1$  and  $x_2$ .

Setting  $x_2 = 1$  (arbitrarily) gives  $x_1 = -2$  and hence

$$\mathbf{x} = \begin{bmatrix} -2 \\ 1 \end{bmatrix}$$

is an eigenvector corresponding to the eigenvalue 1. From the second expression,

$$x_1 = 3x_2$$

$$x_1 + 3x_2 = 3x_2$$

implies  $x_1 = 0$  and  $x_2 = 1$  (arbitrarily), and hence

$$\mathbf{x} = \begin{bmatrix} 0 \\ 1 \end{bmatrix}$$

is an eigenvector corresponding to the eigenvalue 3. It is usual practice to determine an eigenvector so that it has length one. That is, if  $\mathbf{Ax} = \lambda\mathbf{x}$ , we take  $\mathbf{e} = \mathbf{x}/\sqrt{\mathbf{x}'\mathbf{x}}$  as the eigenvector corresponding to  $\lambda$ . For example, the eigenvector for  $\lambda = 1$  is  $[-2/\sqrt{5}, 1/\sqrt{5}]'$ .

**Definition 2A.32.** A quadratic form,  $Q(\mathbf{x})$ , in the  $k$  variables  $x_1, x_2, \dots, x_k$  is  $Q(\mathbf{x}) = \mathbf{x}'\mathbf{A}\mathbf{x}$  where  $\mathbf{x}' = [x_1, x_2, \dots, x_k]$  and  $\mathbf{A}$  is a  $k \times k$  symmetric matrix.

Note that a quadratic form can be written as  $Q(\mathbf{x}) = \sum_{i=1}^k \sum_{j=1}^k a_{ij}x_ix_j$ . For example:

$$Q(\mathbf{x}) = \begin{bmatrix} x_1 & x_2 \end{bmatrix} \begin{bmatrix} 1 & 1 \\ 1 & 1 \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \end{bmatrix} = x_1^2 + 2x_1x_2 + x_2^2$$

$$Q(\mathbf{x}) = \begin{bmatrix} x_1 & x_2 & x_3 \end{bmatrix} \begin{bmatrix} 1 & 3 & 0 \\ 3 & -1 & -2 \\ 0 & -2 & 2 \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \\ x_3 \end{bmatrix} = x_1^2 + 6x_1x_2 - x_2^2 - 4x_2x_3 + 2x_3^2$$

## EXERCISES

2.1. Let  $\mathbf{x}' = [5, 1, 3]$  and  $\mathbf{y}' = [-1, 3, 1]$ .

(a) Graph the two vectors.

(b) Find (i) the length of  $\mathbf{x}$ , (ii) the angle between  $\mathbf{x}$  and  $\mathbf{y}$ , and (iii) the projection of  $\mathbf{y}$  on  $\mathbf{x}$ .

(c) Since  $\bar{x} = 3$  and  $\bar{y} = 1$ , graph  $[5 - 3, 1 - 3, 3 - 3] = [2, -2, 0]$  and  $[-1 - 1, 3 - 1, 1 - 1] = [-2, 2, 0]$ .

2.2. Given the matrices

$$\mathbf{A} = \begin{bmatrix} -1 & 3 \\ 4 & 2 \end{bmatrix}, \quad \mathbf{B} = \begin{bmatrix} 4 & -3 \\ 1 & -2 \\ -2 & 0 \end{bmatrix}, \quad \text{and} \quad \mathbf{C} = \begin{bmatrix} 5 \\ -4 \\ 2 \end{bmatrix}$$

perform the indicated multiplications.

(a)  $5\mathbf{A}$

(b)  $\mathbf{BA}$

(c)  $\mathbf{A}'\mathbf{B}'$

(d)  $\mathbf{C}'\mathbf{B}$

(e) Is  $\mathbf{AB}$  defined?

2.3. Verify the following properties of transpose when

$$\mathbf{A} = \begin{bmatrix} 2 & 1 \\ 1 & 3 \end{bmatrix}, \quad \mathbf{B} = \begin{bmatrix} 1 & 4 & 2 \\ 5 & 0 & 3 \end{bmatrix}, \quad \text{and} \quad \mathbf{C} = \begin{bmatrix} 1 & 4 \\ 3 & 2 \end{bmatrix}$$

(a)  $(\mathbf{A}')' = \mathbf{A}$

(b)  $(\mathbf{C}')^{-1} = (\mathbf{C}^{-1})'$

(c)  $(\mathbf{AB})' = \mathbf{B}'\mathbf{A}'$

(d) Prove (c) for general  $\mathbf{A}$  and  $\mathbf{B}$ .

2.4. When  $\mathbf{A}^{-1}$  and  $\mathbf{B}^{-1}$  exist, prove each of the following.

(a)  $(\mathbf{A}')^{-1} = (\mathbf{A}^{-1})'$

(b)  $(\mathbf{AB})^{-1} = \mathbf{B}^{-1}\mathbf{A}^{-1}$

(Hint: Part a can be proved by noting that  $\mathbf{AA}^{-1} = \mathbf{I}$ ,  $\mathbf{I} = \mathbf{I}'$ , and  $(\mathbf{AA}^{-1})' = (\mathbf{A}^{-1})'\mathbf{A}'$ . Part b follows from  $(\mathbf{B}^{-1}\mathbf{A}^{-1})\mathbf{AB} = \mathbf{B}^{-1}(\mathbf{A}^{-1}\mathbf{A})\mathbf{B} = \mathbf{B}^{-1}\mathbf{B} = \mathbf{I}$ .)

2.5. Check that

$$\mathbf{Q} = \begin{bmatrix} \frac{5}{13} & \frac{12}{13} \\ -\frac{12}{13} & \frac{5}{13} \end{bmatrix}$$

is an orthogonal matrix.

2.6. Let

$$\mathbf{A} = \begin{bmatrix} 9 & -2 \\ -2 & 6 \end{bmatrix}$$

(a) Is  $\mathbf{A}$  symmetric?

(b) Show that  $\mathbf{A}$  is positive definite.

2.7. Let  $\mathbf{A}$  be as given in Exercise 2.6.

(a) Determine the eigenvalues and eigenvectors of  $\mathbf{A}$ .

(b) Write the spectral decomposition of  $\mathbf{A}$ .

(c) Find  $\mathbf{A}^{-1}$ .

(d) Find the eigenvalues and eigenvectors of  $\mathbf{A}^{-1}$ .

2.8. Given the matrix

$$\mathbf{A} = \begin{bmatrix} 1 & 2 \\ 2 & -2 \end{bmatrix}$$

find the eigenvalues  $\lambda_1$  and  $\lambda_2$  and the associated normalized eigenvectors  $\mathbf{e}_1$  and  $\mathbf{e}_2$ . Determine the spectral decomposition (2-16) of  $\mathbf{A}$ .