

1 Principal Component Analysis

1.1 Introduction

Suppose that, as usual, we begin with a collection of measurements of different features for a group of samples. Some of these measurements will tell us quite a bit about the difference among our samples, while others may contain relatively little information. For example, if we are analyzing the effect of a certain weight loss regimen on a group of people, the age and weight of the subjects may have a great deal of influence on how successful the regimen is, while their blood pressure might not. One way to help identify which features are more significant is to ask whether or not the feature varies a lot among the different samples. If nearly all the measurements of a feature are the same, it can't have much power in distinguishing the samples, while if the measurements vary a great deal then that feature has a chance to contain useful information.

In this section we will discuss a way to measure the variability of measurements and then introduce principal component analysis (PCA). PCA is a method for finding which linear combinations of measurements have the greatest variability and therefore might contain the most information. It also allows us to identify combinations of measurements that don't vary much at all. Combining this information, we can sometimes replace our original system of features with a smaller set that still captures most of the interesting information in our data, and thereby find hidden characteristics of the data and simplify our analysis a great deal.

1.2 Variance and Covariance

1.2.1 Variance

Suppose that we have a collection of measurements (x_1, \dots, x_n) of a particular feature X . For example, x_i might be the initial weight of the i th participant in our weight loss study. The mean of the values (x_1, \dots, x_n) is

$$\mu_X = \frac{1}{n} \sum_{i=1}^n x_i.$$

The simplest measure of the variability of the data is called its *variance*.

Definition: The (sample) variance of the data x_1, \dots, x_n is

$$\sigma_X^2 = \frac{1}{n} \sum_{i=1}^n (x_i - \mu_X)^2 = \frac{1}{n} \left(\sum_{i=1}^n x_i^2 \right) - \mu_X^2 \quad (1)$$

The square root of the variance is called the *standard deviation*.

As we see from the formula, the variance is a measure of how ‘spread out’ the data is from the mean.

Recall that in our discussion of linear regression we thought of our set of measurements x_1, \dots, x_n as a vector – it’s one of the columns of our data matrix. From that point of view, the variance has a geometric interpretation – it is $\frac{1}{N}$ times the square of the distance from the point $X = (x_1, \dots, x_n)$ to the point $\mu_X(1, 1, \dots, 1) = \mu_X E$:

$$\sigma_X^2 = \frac{1}{n}(X - \mu_X E) \cdot (X - \mu_X E) = \frac{1}{n} \|X - \mu_X E\|^2. \quad (2)$$

1.2.2 Covariance

The variance measures the dispersion of measures of a single feature. Often, we have measurements of multiple features and we might want to know something about how two features are related. The *covariance* is a measure of whether two features tend to be related, in the sense that when one increases, the other one increases; or when one increases, the other one decreases.

Definition: Given measurements (x_1, \dots, x_n) and (y_1, \dots, y_n) of two features X and Y , the covariance of X and Y is

$$\sigma_{XY} = \frac{1}{N} \sum_{i=1}^N (x_i - \mu_X)(y_i - \mu_Y) \quad (3)$$

There is a nice geometric interpretation of this, as well, in terms of the dot product. If $X = (x_1, \dots, x_n)$ and $Y = (y_1, \dots, y_n)$ then

$$\sigma_{XY} = \frac{1}{N} ((X - \mu_X E) \cdot (Y - \mu_Y E)).$$

From this point of view, we can see that σ_{XY} is positive if the $X - \mu_X E$ and $Y - \mu_Y E$ vectors “point roughly in the same direction” and its negative if they “point roughly in the opposite direction.”

1.2.3 Correlation

One problem with interpreting the variance and covariance is that we don’t have a scale – for example, if σ_{XY} is large and positive, then we’d like to say that X and Y are closely related, but it could be just that the entries of $X - \mu_X E$ and $Y - \mu_Y E$ are large. Here, though, we can really take advantage of the geometric interpretation. Recall that the dot product of two vectors satisfies the formula

$$a \cdot b = \|a\| \|b\| \cos(\theta)$$

where θ is the angle between a and b . So

$$\cos(\theta) = \frac{a \cdot b}{\|a\| \|b\|}.$$

Let's apply this to the variance and covariance, by noticing that

$$\frac{(X - \mu_X E) \cdot (Y - \mu_Y E)}{\|(X - \mu_X E)\| \|(Y - \mu_Y E)\|} = \frac{\sigma_{XY}}{\sigma_X \sigma_Y}$$

so the quantity

$$r_{XY} = \frac{\sigma_{XY}}{\sigma_X \sigma_Y} \tag{4}$$

measures the cosine of the angle between the vectors $X - \mu_X E$ and $Y - \mu_Y E$.

Definition: The quantity r_{XY} defined in eq. 4 is called the (sample) *correlation coefficient* between X and Y . We have $0 \leq |r_{XY}| \leq 1$ with $r_{XY} = \pm 1$ if and only if the two vectors $X - \mu_X E$ and $Y - \mu_Y E$ are collinear in \mathbf{R}^n .

Figure 1 illustrates data with different values of the correlation coefficient.

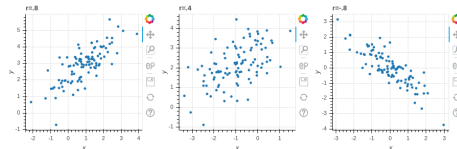


Figure 1: Correlation

1.2.4 The covariance matrix

In a typical situation we have many features for each of our (many) samples, that we organize into a data matrix X . To recall, each column of X corresponds to a feature that we measure, and each row corresponds to a sample. For example, each row of our matrix might correspond to a person enrolled in a study, and the columns correspond to height (cm), weight (kg), systolic blood pressure, and age (in years):

Table 1: A sample data matrix X

sample	Ht	Wgt	Bp	Age
A	180	75	110	35
B	193	80	130	40
...
U	150	92	105	55

If we have multiple features, as in this example, we might be interested in the variance of each feature and all of their mutual covariances. This “package” of information can be obtained “all at once” by taking advantage of some matrix algebra.

Definition: Let X be a $N \times k$ data matrix, where the k columns of X correspond to different features and the N rows to different samples. Let X_0 be the centered version of this data matrix, obtained by subtracting the mean μ_i of column i from all the entries x_{si} in that column. Then the $k \times k$ symmetric matrix

$$D_0 = \frac{1}{N} X_0^\top X_0$$

is called the (sample) covariance matrix for the data.

Proposition: The diagonal entries d_{ii} of D_0 are the variances of the columns of X :

$$d_{ii} = \sigma_i^2 = \frac{1}{N} \sum_{s=1}^N (x_{si} - \mu_i)^2$$

and the off-diagonal entries $d_{ij} = d_{ji}$ are the covariances of the i^{th} and j^{th} columns of X :

$$d_{ij} = \sigma_{ij} = \frac{1}{N} \sum_{s=1}^N (x_{si} - \mu_i)(x_{sj} - \mu_j)$$

The sum of the diagonal entries, the trace of D_0 is the **total** variance of the data.

Proof: This follows from the definitions, but it’s worth checking the details, which we leave as an exercise.

1.2.5 Visualizing the covariance matrix

If the number of features in the data is not too large, a density matrix plot provides a tool for visualizing the covariance matrix of the data. A density matrix plot is an $k \times k$ grid of plots (where k is the number of features). The entry with (i, j) coordinates in the grid is a scatter plot of the i^{th} feature against the j^{th} one if $i \neq j$, and is a histogram of the i^{th} variable if $i = j$.

Figure 2 is an example of a density matrix plot for a dataset with 50 samples and 2 features. This data has been centered, so it can be represented in a 50×2 data matrix X_0 . The upper left and lower right graphs are scatter plots of the two columns, while the lower left and upper right are the histograms of the columns.

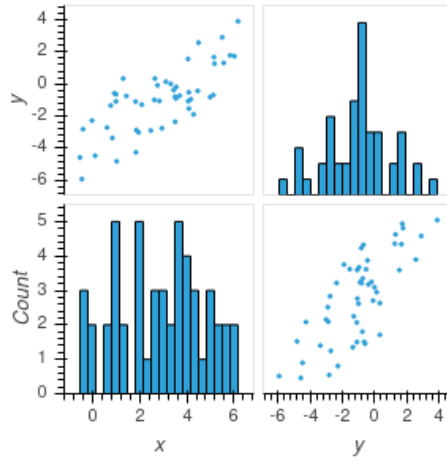


Figure 2: Density Matrix Plot

1.2.6 Linear Combinations of Features (Scores)

Sometimes useful information about our data can be revealed if we combine different measurements together to obtain a “hybrid” measure that captures something interesting. For example, in the Auto MPG dataset that we studied in the section on Linear Regression, we looked at the influence of both vehicle weight w and engine displacement e on gas mileage; perhaps there is some value in considering a hybrid “score” defined as

$$S = aw + be$$

for some constants a and b – maybe by choosing a good combination we could find a better predictor of gas mileage than using one or the other of the features individually.

As another example, suppose we are interested in the impact of the nutritional content of food on weight gain in a study. We know that both calorie content and the level dietary fiber contribute to the weight gain of participants eating this particular food; maybe there is some kind of combined “calorie/fiber” score we could introduce that captures the impact of that food better.

Definition: Let X_0 be a (centered) $N \times k$ data matrix giving information about k features for each of N samples. A linear synthetic feature, or a linear score, is a linear combination of the k features. The linear score is defined by constants a_1, \dots, a_k so that if y_1, \dots, y_k are the values of the features for a particular sample, then the linear score for that sample is

$$S = a_1y_1 + a_2y_2 + \dots + a_ky_k$$

Lemma: The values of the linear score for each of the N samples can be calculated as

$$\begin{bmatrix} S_1 \\ \vdots \\ S_N \end{bmatrix} = X_0 \begin{bmatrix} a_1 \\ \vdots \\ a_k \end{bmatrix}. \quad (5)$$

Proof: Multiplying a matrix by a column vector computes a linear combination of the columns – that’s what this lemma says. Exercise 3 asks you to write out the indices and make sure you believe this.

1.2.7 Mean and variance of scores

When we combine features to make a hybrid score, we assume that the features were centered to begin with, so that each features has mean zero. As a result, the mean of the hybrid features is again zero.

Lemma: A linear combination of features with mean zero again has mean zero.

Proof: Let S_i be the score for the i^{th} sample, so

$$S_i = \sum_{j=1}^k x_{ij} a_j.$$

where X_0 has entries x_{ij} . Then the mean value of the score is

$$\mu_S = \frac{1}{N} \sum_{i=1}^N S_i = \frac{1}{N} \sum_{i=1}^N \sum_{j=1}^k x_{ij} a_j.$$

Reversing the order of the sum yields

$$\mu_S = \frac{1}{N} \sum_{j=1}^k \sum_{i=1}^N x_{ij} a_j = \sum_{j=1}^k a_j \frac{1}{N} \left(\sum_{i=1}^N x_{ij} \right) = \sum_{j=1}^k a_j \mu_j = 0$$

where $\mu_j = 0$ is the mean of the j^{th} feature (column) of X_0 .

The variance is more interesting, and gives us an opportunity to put the covariance matrix to work. Remember from 2 that, since a score S has mean zero, it’s variance is $\sigma_S^2 = \frac{1}{N} S \cdot S$ – where here the score S is represented by the column vector with entries S_1, \dots, S_k as in eq. 5.

Lemma: The variance of the score S with weights a_1, \dots, a_k is

$$\sigma_S^2 = a^\top D_0 a = \begin{bmatrix} a_1 & \cdots & a_k \end{bmatrix} D_0 \begin{bmatrix} a_1 \\ \vdots \\ a_k \end{bmatrix} \quad (6)$$

More generally, if S_1 and S_2 are scores with weights a_1, \dots, a_k and b_1, \dots, b_k respectively, then the covariance $\sigma_{S_1 S_2}$ is

$$\sigma_{S_1 S_2} = a^\top D_0 b.$$

Proof: From eq. 2 and 5 we know that

$$\sigma_S^2 = \frac{1}{N} S \cdot S$$

and

$$S = X_0 a.$$

Since $\frac{1}{N} S \cdot S = \frac{1}{N} S^\top S$, this gives us

$$\frac{1}{N} \sigma_S^2 = \frac{1}{N} (X_0 a)^\top (X_0 a) = \frac{1}{N} a^\top X_0^\top X_0 a = a^\top D_0 a$$

as claimed.

For the covariance, use a similar argument with eq. 3 and eq. 5. writing $\sigma_{S_1 S_2} = \frac{1}{N} S_1 \cdot S_2$ and the fact that S_1 and S_2 can be written as $X_0 a$ and $X_0 b$.

The point of this lemma is that the covariance matrix contains not just the variances and covariances of the original features, but also enough information to construct the variances and covariances for *any linear combination of features*.

In the next section we will see how to exploit this idea to reveal hidden structure in our data.

1.2.8 Geometry of Scores

Let's return to the dataset that we looked at in section 1.2.5. We simplify the density matrix plot in fig. 3, which shows one of the scatter plots and the two histograms.

The scatter plot shows that the data points are arranged in a more or less elliptical cloud oriented at an angle to the xy -axes which represent the two given features. The two individual histograms show the distribution of the two features – each has mean zero, with the x -features distributed between -2 and 2 and the y feature between -4 and 4 . Looking just at the two features individually, meaning only at the two histograms, we can't see the overall elliptical structure.

How can we get a better grip on our data in this situation? We can try to find a “direction” in our data that better illuminates the variation of the data. For example, suppose that we pick a unit vector at the origin pointing in a particular direction in our data. See fig. 4.

Now we can orthogonally project the datapoints onto the line defined by this vector, as shown in fig. 5.

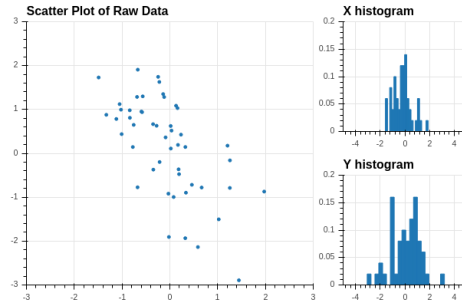


Figure 3: Simulated Data with Two Features

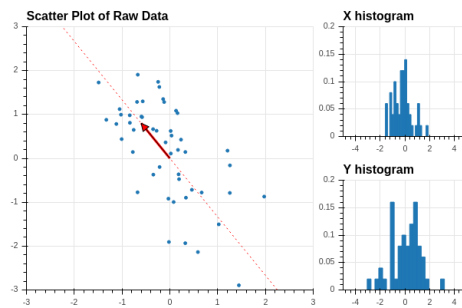


Figure 4: A direction in the data

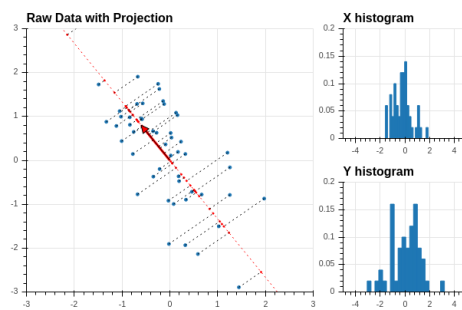


Figure 5: Projecting the datapoints

Recall that if the unit vector is defined by coordinates $u = [u_0, u_1]$, then the orthogonal projection of the point x with coordinates (x_0, x_1) is $(x \cdot u)u$. Now

$$x \cdot u = u_0x_0 + u_1x_1$$

so the coordinates of the points along the line defined by u are the values of the score Z defined by $u = [u_0, u_1]$. Using our work in the previous section, we see that we can find all of these coordinates by matrix multiplication:

$$Z = X_0u$$

where X_0 is our data matrix. Now let's add a histogram of the values of Z to our picture:

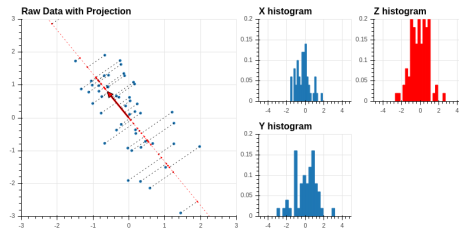


Figure 6: Distribution of Z

This histogram shows the distribution of the values of Z along the tilted line defined by the unit vector u .

Finally, using our work on the covariance matrix, we see that the variance of Z is given by

$$\sigma_Z^2 = \frac{1}{50}u^\top X_0^\top X_0 u = u^\top D_0 u$$

where D_0 is the covariance matrix of the data X_0 .

Lemma: Let X_0 be a $N \times k$ centered data matrix, and let $D_0 = \frac{1}{N}X_0^\top X_0$ be the associated covariance matrix. Let u be a unit vector in “feature space” \mathbf{R}^k . Then the score $S = X_0u$ can be interpreted as the coordinates of the points of X_0 projected onto the line generated by u . The variance of this score is

$$\sigma_S^2 = u^\top D_0 u = \sum_{i=1}^N s_i^2$$

where $s_i = X_0[i, :]u$ is the dot product of the i^{th} row $X_0[i, :]$ with u . It measures the variability in the data “in the direction of the unit vector u .”

1.3 Principal Components

1.3.1 Change of variance with direction

As we’ve seen in the previous section, if we choose a unit vector u in the feature space and find the projection X_0u of our data onto the line through u , we get a

“score” that we can use to measure the variance of the data in the direction of u . What happens as we vary u ?

To study this question, let’s continue with our simulated data from the previous section, and introduce a unit vector

$$u(\theta) = [\cos(\theta) \quad \sin(\theta)].$$

This is in fact a unit vector, since $\sin^2(\theta) + \cos^2(\theta) = 1$, and it is oriented at an angle θ from the x -axis.

The variance of the data in the direction of $u(\theta)$ is given by

$$\sigma_\theta^2 = u(\theta)^\top D_0 u(\theta).$$

A plot of this function for the data we have been considering is in fig. 7. As you can see, the variance goes through two full periods with the angle, and it reaches a maximum and minimum value at intervals of $\pi/2$ – so the two angles where the variance are maximum and minimum are orthogonal to one another.

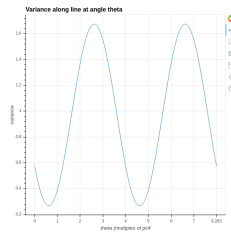


Figure 7: Change of variance with angle theta

The two directions where the variance is maximum and minimum are drawn on the original data scatter plot in fig. 8 .

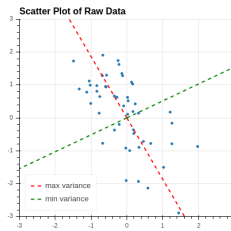


Figure 8: Data with principal directions

Let’s try to understand why this is happening.

1.3.2 Directions of extremal variance

Given our centered, $N \times i$ data matrix X_0 , with its associated covariance matrix $D_0 = \frac{1}{N} X_0^\top X_0$, we would like to find unit vectors u in \mathbf{R}^k so that

$$\sigma_u^2 = u^\top D_0 u$$

reaches its maximum and its minimum. Here σ_u^2 is the variance of the “linear score” $X_0 u$ and it represents how dispersed the data is in the “ u direction” in \mathbf{R}^k .

In this problem, remember that the coordinates of $u = (u_1, \dots, u_k)$ are the variables and the symmetric matrix D_0 is given. As usual, we to find the maximum and minimum values of σ_u^2 , we should look at the partial derivatives of σ_u^2 with respect to the variables u_i and set them to zero. Here, however, there is a catch – we want to restrict u to being a unit vector, with $u \cdot u = \sum u_i^2 = 1$.

So this is a *constrained optimization problem*:

- Find extreme values of the function

$$\sigma_u^2 = u^\top D_0 u$$

- Subject to the constraint $\|u\|^2 = u \cdot u = 1$ (or $u \cdot u - 1 = 0$)

We will use the technique of *Lagrange Multipliers* to solve such a problem.

To apply this method, we introduce the function

$$S(u, \lambda) = u^\top D_0 u - \lambda(u \cdot u - 1) \tag{7}$$

Then we compute the gradient

$$\nabla S = \begin{bmatrix} \frac{\partial S}{\partial u_1} \\ \vdots \\ \frac{\partial S}{\partial u_k} \\ \frac{\partial S}{\partial \lambda} \end{bmatrix} \tag{8}$$

and solve the system of equations $\nabla S = 0$. Here we have written the gradient as a column vector for reasons that will become clearer shortly.

Computing all of these partial derivatives looks messy, but actually if we take advantage of matrix algebra it’s not too bad. The following two lemmas explain how to do this.

Lemma: Let M be a $N \times k$ matrix with constant coefficients and let u be a $k \times 1$ column vector whose entries are u_1, \dots, u_k . The function $F(u) = Mu$ is a

linear map from $\mathbf{R}^k \rightarrow \mathbf{R}^N$. Its (total) derivative is a linear map between the same vector spaces, and satisfies

$$D(F)(v) = Mv$$

for any $k \times 1$ vector v . If u is a $1 \times N$ matrix, and $G(u) = uM$, then

$$D(G)(v) = vM$$

for any $1 \times N$ vector v . (This is the matrix version of the derivative rule that $\frac{d}{dx}(ax) = a$ for a constant a .)

Proof: Since $F : \mathbf{R}^k \rightarrow \mathbf{R}^N$, we can write out F in more traditional function notation as

$$F(u) = (F_1(u_1, \dots, u_k), \dots, F_N(u_1, \dots, u_k))$$

where

$$F_i(u_1, \dots, u_k) = \sum_{j=1}^k m_{ij}u_j.$$

Thus $\frac{\partial F_i}{\partial u_j} = m_{ij}$. The total derivative $D(F)$ is the linear map with matrix

$$D(F)_{ij} = \frac{\partial F_i}{\partial u_j} = m_{ij}$$

and so $D(F) = M$.

The other result is proved the same way.

Lemma: Let D be a symmetric $k \times k$ matrix with constant entries and let u be an $k \times 1$ column vector of variables u_1, \dots, u_k . Let $F : \mathbf{R}^k \rightarrow R$ be the function $F(u) = u^T D u$. Then the gradient $\nabla_u F$ is a vector field – that is, a vector-valued function of u , and is given by the formula

$$\nabla_u F = 2Du$$

Proof: Let d_{ij} be the i, j entry of D . We can write out the function F to obtain

$$F(u_1, \dots, u_k) = \sum_{i=1}^k \sum_{j=1}^k u_i d_{ij} u_j.$$

Now $\frac{\partial F}{\partial u_i}$ is going to pick out only terms where u_i appears, yielding:

$$\frac{\partial F}{\partial u_i} = \sum_{j=1}^k d_{ij} u_j + \sum_{j=1}^k u_j d_{ji}$$

Here the first sum catches all of the terms where the first “ u ” is u_i ; and the second sum catches all the terms where the second “ u ” is u_i . The diagonal terms

$u_i^2 d_{ii}$ contribute once to each sum, which is consistent with the rule that the derivative of $u_i^2 d_{ii} = 2u_i d_{ii}$. To finish the proof, notice that

$$\sum_{j=1}^k u_j d_{ji} = \sum_{j=1}^k d_{ij} u_j$$

since D is symmetric, so in fact the two terms are the same. Thus

$$\frac{\partial}{\partial u_i} F = 2 \sum_{j=1}^k d_{ij} u_j$$

But the right hand side of this equation is twice the i^{th} entry of Du , so putting the results together we get

$$\nabla_u F = \begin{bmatrix} \frac{\partial F}{\partial u_1} \\ \vdots \\ \frac{\partial F}{\partial u_k} \end{bmatrix} = 2Du.$$

The following theorem puts all of this work together to reduce our questions about how variance changes with direction.

1.3.3 Critical values of the variance

Theorem: The critical values of the variance σ_u^2 , as u varies over unit vectors in \mathbf{R}^N , are the eigenvalues $\lambda_1, \dots, \lambda_k$ of the covariance matrix D , and if e_i is a unit eigenvector corresponding to λ_i , then $\sigma_{e_i}^2 = \lambda_i$.

Proof: Recall that we introduced the Lagrange function $S(u, \lambda)$, whose critical points give us the solutions to our constrained optimization problem. As we said in eq. 7:

$$S(u, \lambda) = u^T D_0 u - \lambda(u \cdot u - 1) = u^T D_0 u - \lambda(u \cdot u) + \lambda$$

Now apply our Matrix calculus lemmas. First, let's treat λ as a constant and focus on the u variables. We can write $u \cdot u = u^T I_N u$ where I_N is the identity matrix to compute:

$$\nabla_u S = 2D_0 u - 2\lambda u$$

For λ we have

$$\frac{\partial}{\partial \lambda} S = -u \cdot u + 1.$$

The critical points occur when

$$\nabla_u S = 2(D_0 - \lambda)u = 0$$

and

$$\frac{\partial}{\partial \lambda} S = 1 - u \cdot u = 0$$

The first equation says that λ must be an eigenvalue, and u an eigenvector:

$$D_0 u = \lambda u$$

while the second says u must be a unit vector $u \cdot u = \|u\|^2 = 1$. The second part of the result follows from the fact that if e_i is a unit eigenvector with eigenvalue λ_i then

$$\sigma_{e_i}^2 = e_i^\top D_0 e_i = \lambda_i \|e_i\|^2 = \lambda_i.$$

To really make this result pay off, we need to recall some key facts about the eigenvalues and eigenvectors of symmetric matrices. Because these facts are so central to this result, and to other applications throughout machine learning and mathematics generally, we provide proofs in section 1.5.

Table 2: Properties of Eigenvalues of Real Symmetric Matrices

Summary

1. All of the eigenvalues $\lambda_1, \dots, \lambda_l$ of D are real. If $u^\top D u \geq 0$ for all $u \in \mathbf{R}^k$, then all eigenvalues λ_i are non-negative. In the latter case we say that D is *positive semi-definite*.
 2. If v is an eigenvector for D with eigenvalue λ , and w is an eigenvector with a different eigenvalue λ' , then v and w are orthogonal: $v \cdot w = 0$.
 3. There is an orthonormal basis u_1, \dots, u_k of \mathbf{R}^k made up of eigenvectors of D corresponding to the eigenvalues λ_i .
 4. Let Λ be the diagonal matrix with entries $\lambda_1, \dots, \lambda_N$ and let P be the matrix whose columns are made up of the vectors u_i . Then $D = P \Lambda P^\top$.
-

If we combine this theorem with the facts summarized in table 2 then we get a complete picture. Let D_0 be the covariance matrix of our data. Since

$$\sigma_u^2 = u^\top D_0 u \geq 0 \text{ (it's a sum of squares)}$$

we know that the eigenvalues $\lambda_1 \geq \lambda_2 \geq \dots \geq \lambda_k \geq 0$ are all nonnegative. Choose a corresponding sequence u_1, \dots, u_k of orthogonal eigenvectors where all $\|u_i\|^2 = 1$. Since the u_i form a basis of \mathbf{R}^N , any score is a linear combination of the u_i :

$$S = \sum_{i=1}^k a_i u_i.$$

Since $u_i^\top D_0 u_j = \lambda_j u_i^\top u_j = 0$ unless $i = j$, in which case it is λ_i , we can compute

$$\sigma_S^2 = \sum_{i=1}^k \lambda_i a_i^2,$$

and $\|S\|^2 = \sum_{i=1}^k a_i^2$ since the u_i are an orthonormal set. So in these coordinates, our optimization problem is:

- maximize $\sum \lambda_i a_i^2$
- subject to the constraint $\sum a_i^2 = 1$.

We don't need any fancy math to see that the maximum happens when $a_1 = 1$ and the other $a_j = 0$, and in that case, the maximum is λ_1 . (If λ_1 occurs more than once, there may be a whole subspace of directions where the variance is maximal). Similarly, the minimum value is λ_k and occurs when $a_k = 1$ and the others are zero.

1.3.4 Subspaces of extremal variance

We can generalize the question asked in section 1.3.2 by seeking, not just a vector u pointing in the direction of the extremal variance, but instead the *subspace* U_s of dimension s with the property that the total variance of the projection of the data into U_s is maximal compared to its projection into other subspaces of that dimension.

To make this concrete, suppose we consider a subspace E of \mathbf{R}^k of dimension t with basis w_1, \dots, w_t . Complete this to a basis $w_1, \dots, w_t, w_{t+1}, \dots, w_k$ of \mathbf{R}^k and then apply the Gram Schmidt Process (see section 1.5.1) to find an orthonormal basis $w'_1, \dots, w'_s, w'_{s+1}, \dots, w'_k$ where the w'_1, \dots, w'_t are an orthonormal basis for E . Let W be the $k \times t$ matrix whose columns are the w'_i for $i = 1, \dots, t$. The rows of the matrix $X_0 W$ give the coordinates of the projection of each sample into the subspace E expressed in terms of the scores corresponding to these vectors w'_i . The total variance of these projections is

$$\sigma_E^2 = \sum_{i=1}^t \|X_0 w'_i\|^2 = \sum_{i=1}^t (w'_i)^\top X_0^\top X_0 w'_i = \sum_{i=1}^t (w'_i)^\top D_0 w'_i$$

If we want to maximize this, we have the constrained optimization problem of finding w'_1, \dots, w'_t so that

- $\sum_{i=1}^t (w'_i)^\top D_0 w'_i$ is maximal
- subject to the constraint that each w_i has $\|w'_i\|^2 = 1$,
- and that the w'_i are orthogonal, meaning $w'_i \cdot w'_j = 0$ for $i \neq j$,
- and that the w'_i are linearly independent.

Then the span E of these w'_i is subspace of extremal variance.

Theorem: A t -dimensional subspace E is a subspace of extremal variance if and only if it is spanned by t orthonormal eigenvectors of the matrix D_0 corresponding to the t largest eigenvalues for D_0 .

Proof: We can approach this problem using Lagrange multipliers and matrix calculus if we are careful. Our unknown is $k \times t$ matrix W whose columns are the

t (unknown) vectors w'_i . The objective function that we are seeking to maximize is

$$F = \text{trace}(W^\top D_0 W) = \sum_{i=1}^t (w'_i)^\top D_0 w_i.$$

The constraints are the requirements that $\|w'_i\|^2 = 1$ and $w'_i \cdot w'_j = 0$ if $i \neq j$. If we introduce a matrix of Lagrange multipliers $\Lambda = (\lambda_{ij})$, where λ_{ij} is the multiplier that goes with the first of these constraints when $i = j$, and the second when $i \neq j$, we can express our Lagrange function as:

$$S(W, \Lambda) = \text{trace}(W^\top D_0 W) - (W^\top W - I)\Lambda$$

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

Taking the derivatives with respect to the entries of W and of Λ yields the following two equations:

$$\begin{aligned} D_0 W &= W\Lambda \\ W^\top W &= I \end{aligned}$$

The first of these equations says that the space E spanned by the columns of W is *invariant* under D_0 , while the second says that the columns of W form an orthonormal basis.

Let's assume for the moment that we have a matrix W that satisfies these conditions.

Then it must be the case that Λ is a symmetric, real valued $t \times t$ matrix, since

$$W^\top D_0 W = W^\top W\Lambda = \Lambda.$$

and the matrix on the left is symmetric.

By the properties of real symmetric matrices (the spectral theorem), there are orthonormal vectors q_1, \dots, q_t that are eigenvectors of Λ with corresponding eigenvalues τ_i . If we let Q be the matrix whose columns are the vectors q_i and let T be the diagonal $t \times t$ matrix whose entries are the τ_i , we have

$$\Lambda Q = QT.$$

If we go back to our original equations, we see that if W exists such that $DW = W\Lambda$, then there is a matrix Q with orthonormal columns and a diagonal matrix T such that

$$D_0 W Q = W\Lambda Q = WQT.$$

In other words, WQ is a matrix whose columns are eigenvectors of D_0 with eigenvalues τ_i for $i = 1, \dots, t$.

Thus we see how to construct an invariant subspace E and a solution matrix W . Such an E is spanned by t orthonormal eigenvectors q_i with eigenvalues τ_i of D_0 ; and W is the matrix whose columns are the q_i . Further, in that case, the total variance associated to E is the sum of the eigenvalues τ_i ; to make this as large as possible, we should choose our eigenvectors to correspond to t of the largest eigenvalues of D_0 . This concludes the proof.

1.3.5 Definition of Principal Components

Definition: The orthonormal unit eigenvectors u_i for D_0 are the *principal directions* or *principal components* for the data X_0 .

Theorem: The maximum variance occurs in the principal direction(s) associated to the largest eigenvalue, and the minimum variance in the principal direction(s) associated with the smallest one. The covariance between scores in principal directions associated with different eigenvalues is zero.

At this point, the picture in fig. 8 makes sense – the red and green dashed lines are the principal directions, they are orthogonal to one another, and the point in the directions where the data is most (and least) “spread out.”

Proof: The statement about the largest and smallest eigenvalues is proved at the very end of the last section. The covariance of two scores corresponding to different eigenvectors u_i and u_j is

$$u_i^\top D_0 u_j = \lambda_j (u_i \cdot u_j) = 0$$

since the u_i and u_j are orthogonal.

Sometimes the results above are presented in a slightly different form, and may be referred to, in part, as Rayleigh’s theorem.

Corollary: (Rayleigh’s Theorem) Let D be a real symmetric matrix and let

$$H(v) = \max_{v \neq 0} \frac{v^\top D v}{v^\top v}.$$

Then $H(v)$ is the largest eigenvalue of D . (Similarly, if we replace max by min, then the minimum is the least eigenvalue).

Proof: The maximum of the function $H(v)$ is the solution to the same optimization problem that we considered above.

Exercises.

1. Prove that the two expressions for σ_X^2 given in section 1.2.1 are the same.
2. Prove that the covariance matrix is as described in the proposition in 1.2.4.
3. Let X_0 be a $k \times N$ matrix with entries x_{ij} for $1 \leq i \leq k$ and $1 \leq j \leq N$. If a linear score is defined by the constants a_1, \dots, a_N , check that equation eq. 5 holds as claimed.

4. Why is it important to use a unit vector when computing the variance of X_0 in the direction of u ? Suppose $v = \lambda u$ where u is a unit vector and $\lambda > 0$ is a constant. Let S' be the score $X_0 v$. How is the variance of S' related to that of $S = X_0 u$?

1.4 Dimensionality Reduction via Principal Components

The principal components associated with a dataset separate out directions in the feature space in which the data is most (or least) variable. One of the main applications of this information is to enable us to take data with a great many features – a set of points in a high dimensional space – and, by focusing our attention on the scores corresponding to the principal directions, capture most of the information in the data in a much lower dimensional setting.

To illustrate how this is done, let X be a $N \times k$ data matrix, let X_0 be its centered version, and let $D_0 = \frac{1}{N} X_0^T X$ be the associated covariance matrix.

Apply the spectral theorem (described in table 2) and proved in section 1.5 to the covariance matrix to obtain eigenvalues $\lambda_1 \geq \lambda_2 \geq \dots \lambda_k \geq 0$ and associated eigenvectors u_1, \dots, u_k . The scores $S_i = X_0 u_i$ give the values of the data in the principal directions. The variance of S_i is λ_i .

Now choose a number $t < k$ and consider the vectors S_1, \dots, S_t . The j^{th} entry in S_i is the value of the score S_i for the j^{th} data point. Because S_1, \dots, S_t capture the most significant variability in the original data, we can learn a lot about our data by considering just these t features of the data, instead of needing all N .

To illustrate, let's look at an example. We begin with a synthetic dataset X_0 which has 200 samples and 15 features. The data (some of it) for some of the samples is shown in table 3.

Table 3: Simulated Data for PCA Analysis

	f-0	f-1	f-2	f-3	f-4	...	f-10	f-11	f-12	f-13	f-14
s-0	1.18	-0.41	2.02	0.44	2.24	...	0.32	0.95	0.88	1.10	0.89
s-1	0.74	0.58	1.54	0.23	2.05	...	0.99	1.14	1.56	0.99	0.59
...
s-198	1.04	2.02	1.44	0.40	1.33	...	0.62	0.62	0.54	1.96	0.04
s-199	0.92	2.09	1.58	1.19	1.17	...	0.42	0.85	0.83	2.22	0.90

The full dataset is a 200×15 matrix; it has 3000 numbers in it and we're not really equipped to make sense of it. We could try some graphing – for example, fig. 9 shows a scatter plot of two of the features plotted against each other.

Unfortunately there's not much to see in fig. 9 – just a blob – because the individual features of the data don't tell us much in isolation, whatever structure there is in this data arises out of the relationship between different features.

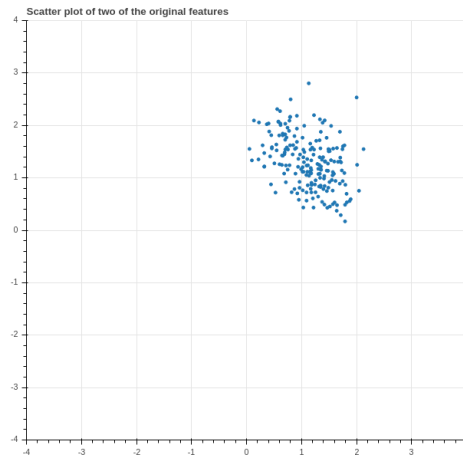


Figure 9: Scatter Plot of Two Features

In fig. 10 we show a “density grid” plot of the data. The graph in position i, j shows a scatter plot of the i^{th} and j^{th} columns of the data, except in the diagonal positions, where in position i, i we plot a histogram of column i . There’s not much structure visible; it is a lot of blobs.

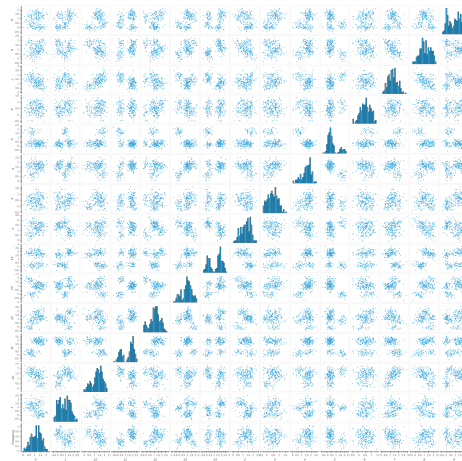


Figure 10: Density Grid Plot of All Features

So let’s apply the theory of principal components. We use a software package to compute the eigenvalues and eigenvectors of the matrix D_0 . The 15 eigenvalues $\lambda_1 \geq \dots \geq \lambda_{15}$ are plotted, in descending order, in fig. 11 .

This plot shows that the first 4 eigenvalues are relatively large, while the remaining 11 are smaller and not much different from each other. We interpret this as

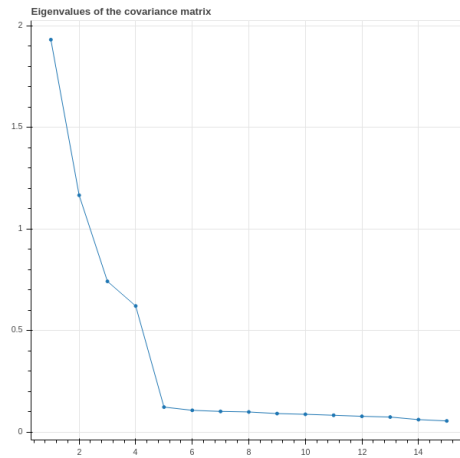


Figure 11: Eigenvalues of the Covariance Matrix

saying that *most of the variation in the data is accounted for by the first four principal components*. We can even make this quantitative. The *total variance* of the data is the sum of the eigenvalues of the covariance matrix – the trace of D_0 – and in this example that sum is around 5. The sum of the first 4 eigenvalues is about 4, so the first four eigenvalues account for about $4/5$ of the total variance, or about 80% of the variation of the data.

Now let's focus in on the two largest eigenvalues λ_1 and λ_2 and their corresponding eigenvectors u_1 and u_2 . The 200×1 column vectors $S_1 = X_0 u_1$ and $S_2 = X_0 u_2$ are the values of the scores associated with these two eigenvectors. So for each data point (each row of X_0) we have two values (the corresponding entries of S_1 and S_2 .) In fig. 12 we show a scatter plot of these scores.

Notice that suddenly some structure emerges in our data! We can see that the 200 points are separated into five clusters, distinguished by the values of their scores! This ability to find hidden structure in complicated data, is one of the most important applications of principal components.

If we were dealing with real data, we would now want to investigate the different groups of points to see if we can understand what characteristics the principal components have identified.

1.4.1 Loadings

There's one last piece of the PCA puzzle that we are going to investigate. In fig. 12, we plotted our data points in the coordinates given by the first two principal components. In geometric terms, we took the cloud of 200 points in \mathbf{R}^{15} given by the rows of X_0 and projected those points into the two dimensional plane spanned by the eigenvectors u_1 and u_2 , and then plotted the distribution

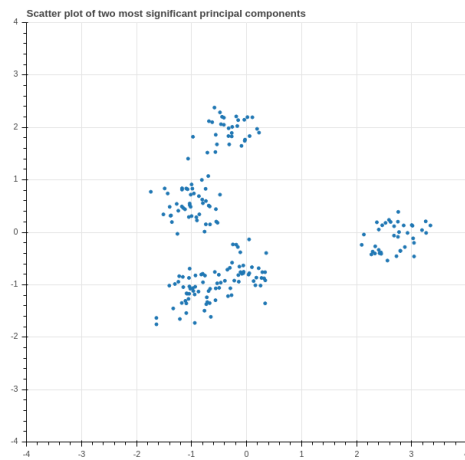


Figure 12: Scatter Plot of Scores in the First Two Principal Directions

of the points in that plane.

More generally, suppose we take our dataset X_0 and consider the first t principal components corresponding to the eigenvectors u_1, \dots, u_t . The projection of the data into the space spanned by these eigenvectors is represented by the $S = k \times t$ matrix X_0U where U is the $k \times t$ matrix whose columns are the eigenvectors u_i . Each row of S gives the values of the score arising from u_i in the i^{th} column for $i = 1, \dots, t$.

The remaining question that we wish to consider is: how can we see some evidence of the original features in subspace? We can answer this by imagining that we had an artificial sample x that has a measurement of 1 for the i^{th} feature and a measurement of zero for all the other features. The corresponding point is represented by a $1 \times k$ row vector with a 1 in position i . The projection of this synthetic sample into the span of the first t principal components is the $1 \times t$ vector xU . Notice, however, that xU is just the i^{th} row of the matrix U . This vector in the space spanned by the u_i is called the “loading” of the i^{th} feature in the principal components.

This is illustrated in section 1.4.1, which shows a line along the direction of the loading corresponding to the each feature added to the scatter plot of the data in the plane spanned by the first two principal components. One observation one can make is that some of the features are more “left to right,” like features 7 and 8, while others are more “top to bottom,” like 6. So points that lie on the left side of the plot have smaller values of features 7 and 8, while those at the top of the plot have larger values of feature 6.

In the next, and last, section, of this discussion of Principal Component Analysis, we will give proofs of the key mathematical ideas summarized earlier in table 2,

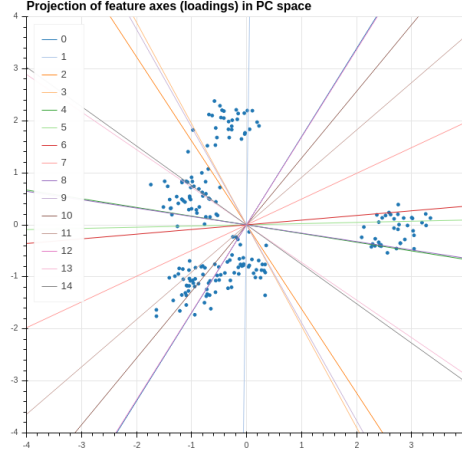


Figure 13: Loadings in the Principal Component Plane

which have been central to this analysis.

1.4.2 The singular value decomposition

The singular value decomposition is a slightly different way of looking at principal components. Let Λ be the diagonal matrix of eigenvalues of D_0 ; we know that the entries of D_0 are non-negative. Let's drop the eigenvectors corresponding to the zero eigenvalue. Let's say that there are s non-zero eigenvalues, and s corresponding eigenvectors.

Lemma: Let P' be the $N \times s$ matrix whose columns are the eigenvectors with non-zero eigenvalues, and let Λ_+ be the $s \times s$ diagonal matrix whose entries are the non-zero eigenvalues. Then $P'\Lambda_+P'^\top = P\Lambda P^\top = D_0$.

Proof: First observe that $P'\Lambda_+P'^\top$ is in fact an $N \times N$ matrix. Then look at the block structure to verify the result.

The matrix Λ_+ is diagonal, invertible, and, since the eigenvalues are positive, it makes sense to consider the real matrix $\Lambda_+^{1/2}$ whose entries are the square roots of the eigenvalues.

Let $U = X_0P'\Lambda_+^{-1/2}$. Note that U is a $k \times s$ dimensional matrix.

Lemma: The columns of U are orthonormal.

Proof: Compute the $s \times s$ matrix $U^\top U$, whose entries are all of the dot products of the columns of U :

$$\begin{aligned} U^\top U &= \Lambda_+^{-1/2} P'^\top X_0^\top X_0 P' \Lambda_+^{-1/2} \\ &= \Lambda_+^{-1/2} P'^\top P' \Lambda_+ P'^\top P' \Lambda_+^{-1/2} \\ &= I_s \end{aligned}$$

by the previous lemma and the fact that $P'P'^\top$ is the $s \times s$ identity matrix.

Rearranging this yields the singular value decomposition.

Theorem: (The singular value decomposition) The matrix X_0 has a decomposition:

$$X_0 = U\Lambda_+^{-1/2}P'^\top$$

where U (of dimension $k \times s$) and P' (of dimension $N \times s$) are orthogonal, and Λ_+ (of dimension $s \times s$) is diagonal with positive entries. Furthermore, the entries of Λ_+ are the non-negative eigenvalues of $D_0 = X_0^\top X_0$, and U and P' are uniquely determined by X_0 .

Proof: We won't work through all of this, as it is a reinterpretation of our work on principal components.

Remark: The entries of $\Lambda_+^{-1/2}$ are called the singular values of X_0 . They can be found directly by considering the optimization problem implicitly equivalent to the problem we solved in section 1.3.4.

1.5 Eigenvalues and Eigenvectors of Real Symmetric Matrices (The Spectral Theorem)

Now that we've shown how to apply the theory of eigenvalues and eigenvectors of symmetric matrices to extract principal directions from data, and to use those principal directions to find structure, we will give a proof of the properties that we summarized in table 2.

A key tool in the proof is the Gram-Schmidt orthogonalization process.

1.5.1 Gram-Schmidt

Proposition (Gram-Schmidt Process): Let w_1, \dots, w_k be a collection of linearly independent vectors in \mathbf{R}^N and let W be the span of the w_i . Let $u_1 = w_1$ and let

$$u_i = w_i - \sum_{j=1}^{i-1} \frac{w_i \cdot u_j}{u_j \cdot u_j} u_j$$

for $i = 2, \dots, k$. Then

- The vectors u_i are orthogonal: $u_i \cdot u_j = 0$ unless $i = j$.
- The vectors u_i span W .
- Each u_i is orthogonal to all of w_1, \dots, w_{i-1} .
- The vectors $u'_i = u_i / \|u_i\|$ are orthonormal.

Proof: This is an inductive exercise, and we leave it to you to work out the details.

1.5.2 The spectral theorem

Theorem: Let D be a real symmetric $N \times N$ matrix. Then:

1. All of the N eigenvalues $\lambda_1 \geq \lambda_2 \geq \dots \geq \lambda_N$ are real. If $u^\top Du \geq 0$ for all $u \in \mathbf{R}^N$, then all eigenvalues $\lambda_i \geq 0$.
2. The matrix D is diagonalizable – that is, it has N linearly independent eigenvectors.
3. If v and w are eigenvectors corresponding to eigenvalues λ and λ' , with $\lambda \neq \lambda'$, then v and w are orthogonal: $v \cdot w = 0$.
4. There is an orthonormal basis u_1, \dots, u_N of \mathbf{R}^N made up of eigenvectors for the eigenvalues λ_i .
5. Let Λ be the diagonal matrix with entries $\lambda_1, \dots, \lambda_N$ and let P be the matrix whose columns are made up of the eigenvectors u_i . Then $D = P\Lambda P^\top$.

Proof: Start with 1. Suppose that λ is an eigenvalue of D . Let u be a corresponding nonzero eigenvector. Then $Du = \lambda u$ and $D\bar{u} = \bar{\lambda}\bar{u}$, where \bar{u} is the vector whose entries are the conjugates of the entries of u (and $\bar{\bar{D}} = D$ since D is real). Now we have

$$\bar{u}^\top Du = \lambda \bar{u} \cdot u = \lambda \|u\|^2$$

and

$$u^\top D\bar{u} = \bar{\lambda} u \cdot \bar{u} = \bar{\lambda} \|u\|^2.$$

But the left hand side of both of these equations are the same (take the transpose and use the symmetry of D) so we must have $\lambda \|u\|^2 = \bar{\lambda} \|u\|^2$ so $\lambda = \bar{\lambda}$, meaning λ is real.

If we have the additional property that $u^\top Du \geq 0$ for all u , then in particular $u_i^\top Du_i = \lambda \|u_i\|^2 \geq 0$, and since $\|u_i\|^2 > 0$ we must have $\lambda \geq 0$.

Property 2 is in some ways the most critical fact. We know from the general theory of the characteristic polynomial, and the fundamental theorem of algebra, that D has N complex eigenvalues, although some may be repeated. However, it may not be the case that D has N linearly independent eigenvectors – it may not be *diagonalizable*. So we will establish that now.

A one-by-one matrix is automatically symmetric and diagonalizable. In the N -dimensional case, we know, at least, that D has at least one eigenvector, and real one at that by part 1, and this gives us a place to begin an inductive argument.

Let $v_N \neq 0$ be an eigenvector with eigenvalue λ and normalized so that $\|v_N\|^2 = 1$,

and extend this to a basis v_1, \dots, v_N of \mathbf{R}^N . Apply the Gram-Schmidt process to construct an orthonormal basis of \mathbf{R}^N u_1, \dots, u_N so that $u_N = v_N$.

Any vector $v \in \mathbf{R}^N$ is a linear combination

$$v = \sum_{i=1}^N a_i u_i$$

and, since the u_i are orthonormal, the coefficients can be calculated as $a_i = (u_i \cdot v)$.

Using this, we can find the matrix D' of the linear map defined by our original matrix D in this new basis. By definition, if d'_{ij} are the entries of D' , then

$$Du_i = \sum_{j=1}^N d'_{ij} u_j$$

and so

$$d'_{ij} = u_j \cdot Du_i = u_j^T Du_i.$$

Since D is symmetric, $u_j^T Du_i = u_i^T Du_j$ and so $d'_{ij} = d'_{ji}$. In other words, the matrix D' is still symmetric. Furthermore,

$$d'_{Ni} = u_i \cdot Du_N = u_i \cdot \lambda u_N = \lambda (u_i \cdot u_N)$$

since $u_N = v_N$. Since the u_i are an orthonormal basis, we see that $d'_{iN} = 0$ unless $i = N$, and $d'_{NN} = \lambda$.

In other words, the matrix D' has a block form:

$$D' = \begin{pmatrix} * & * & \cdots & * & 0 \\ \vdots & \vdots & \ddots & \vdots & \vdots \\ * & * & \cdots & * & 0 \\ 0 & 0 & \cdots & 0 & \lambda \end{pmatrix}$$

and the block denoted by *'s is symmetric. If we call that block D_* , the inductive hypothesis tells us that the symmetric matrix D_* is diagonalizable, so it has a basis of eigenvectors u'_1, \dots, u'_{N-1} with eigenvalues $\lambda_1, \dots, \lambda_{N-1}$; this gives us a basis for the subspace of \mathbf{R}^N spanned by u_1, \dots, u_{N-1} which, together with u_N gives us a basis of \mathbf{R}^N consisting of eigenvectors of D .

This finishes the proof of Property 2.

For property 3, compute

$$v^T Dw = \lambda'(v \cdot w) = w^T Dv = \lambda(w \cdot v).$$

Since $\lambda \neq \lambda'$, we must have $v \cdot w = 0$.

For property 4, if the eigenvalues are all distinct, this is a consequence of property 2 – you have N eigenvectors, scaled to length 1, for different eigenvalues, and by 2 they are orthogonal. So the only complication is the case where some eigenvalues are repeated. If λ occurs r times, then you have r linearly independent vectors u_1, \dots, u_r that span the λ eigenspace. The Gram-Schmidt process allows you to construct an orthonormal set that spans this eigenspace, and while this orthonormal set isn't unique, any one of them will do.

For property 5, let e_i be the column vector that is zero except for a 1 in position i . The product $e_j^\top D e_i = d_{ij}$. Let's write e_i and e_j in terms of the orthonormal basis u_1, \dots, u_N :

$$e_i = \sum_{k=1}^N (e_i \cdot u_k) u_k \text{ and } e_j = \sum_{k=1}^N (e_j \cdot u_k) u_k.$$

Using this expansion, we compute $e_j^\top D e_i$ in a more complicated way:

$$e_j^\top D e_i = \sum_{r=1}^N \sum_{s=1}^N (e_j \cdot u_r) (e_i \cdot u_s) (u_r^\top D u_s).$$

But $u_r^\top D u_s = \lambda_s (u_r \cdot u_s) = 0$ unless $r = s$, in which case it equals λ_r , so

$$e_j^\top D e_i = \sum_{r=1}^N \lambda_r (e_j \cdot u_r) (e_i \cdot u_r).$$

On the other hand,

$$P^\top e_i = \begin{bmatrix} (e_i \cdot u_1) \\ (e_i \cdot u_2) \\ \vdots \\ (e_i \cdot u_N) \end{bmatrix}$$

and

$$\Lambda P^\top e_i = \begin{bmatrix} \lambda_1 (e_i \cdot u_1) \\ \lambda_2 (e_i \cdot u_2) \\ \vdots \\ \lambda_N (e_i \cdot u_N) \end{bmatrix}$$

Therefore the i, j entry of $P \Lambda P^\top$ is

$$(e_j^\top P) \Lambda (P^\top e_i) = \sum_{r=1}^N \lambda_r (e_j \cdot u_r) (e_i \cdot u_r) = d_{ij}$$

so the two matrices D and $P \Lambda P^\top$ are in fact equal.

Exercises:

1. Prove the rest of the first lemma in section 1.4.2.
2. Prove the Gram-Schmidt Process has the claimed properties in section 1.5.1.