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## Principal Components

## Introduction

Principal component analysis is an approach to the problem of finding simple approximate descriptions of variations in large datasets. More concretely, suppose there are $n$ time series of length $T$ :

$$
X_{t, k}, \quad t=1, \ldots, T, \quad k=1, \ldots, n
$$

These could be, among other things, the daily returns on $n$ assets, or the LIBOR rates for $n$ different loan durations. You could try to "explain" these time series in terms of a single series $U_{t}$. This might involve choosing an optimal "weight" $w_{k}$ for series $k$ and minimizing the residuals, $R_{t, k}$ :

$$
\begin{equation*}
X_{t, k}=w_{k} U_{t}+R_{t, k} \tag{1}
\end{equation*}
$$

If you have the "explanation" series $U_{t}$, then you can look for the optimal weight $w_{k}$ for each series. If you have $n$ series, you could look for single series $V_{t}$ that does the best overall. This would be the first principal component for the collection of series $X_{t, k}$.

The time series $U_{t}$ is the single time series that has the most in common with the $n$ separate series $X_{t, k}$ for $k=1, \ldots, n$. If $X_{t, k}$ is the daily return for asset $k$ on day $t$, then $U_{t}$ can be interpreted as the overall market. The weight $w_{k}$ is a multiplier that describes how much "influence" the overall market variable $U_{t}$ has on asset $k$. If $X_{t, k}$ is the LIBOR rate on day $t$ for a loan with duration $D_{k}$ (in days), then $U_{t}$ represents the overall interest rate picture on day $t$. It answers the question: "what do interest rates look like today?" Suppose $D_{1}=1$ for the overnight interest rate and $D_{n}=9958$ for 30 year loans. Then the series $U_{t}$ represents the tendency for overnight and 30 year rates to rise and fall together. The weight $w_{k}$ says how sensitive duration $D_{k}$ loans are to the overall interest rate variable $U_{t}$. If short term interest rates fluctuate less than long rates, then $w_{1}$ will be smaller than $w_{n}$. It could happen (but doesn't) that the long rate goes down when the short rate goes up. In that case $w_{1}$ and $w_{n}$ would have different signs.

You can look for more than one principal component. There could be time series $U_{t, j}$ for $j=1, \ldots, m$ for $m$ principal components. Then time series $X_{t, k}$ would have weight $w_{j, k}$ for the principal components $j=1, \ldots, m$. The residuals would be defined by

$$
X_{t, k}=\sum_{j=1}^{m} w_{k j} U_{t, j}+R_{t, k}
$$

In the interest rate example, $U_{t, 1}$ might represent the overall "level" of interest rates on day $t$ while $U_{t, 2}$ could represent the "slope" of the yield curve. If $U_{t, 2}$ is large then $X_{t, n}$ (long rate) is much larger than $X_{t, 1}$ (short rate).

Principal component analysis means finding and interpreting the principal components. It uses the residual sum of squares optimality criterion. The principal components and weights are chosen to minimize

$$
\begin{equation*}
Q_{m}=\sum_{t=1}^{T} \sum_{k=1}^{n} R_{t, k}^{2} \tag{2}
\end{equation*}
$$

This is done sequentially. The first principal component $U_{t, 1}$ is chosen to minimize $Q_{1}$, then the second principal component $U_{t, 2}$ is chosen to minimize $Q_{2}$, and so on.

The index $t$ need not refer to time. For example, if $k$ refers to a company, then the vector $X_{k}$ could consist of $T$ numbers associated to that company. The numbers could be things like earnings, debt, profit, market capitalization, etc. Then principal component analysis would mean making company "profiles", or looking for the ways the variation between companies can be explained by a small number of factors. Images are another application. Here, image $k$ is described by, say, $T=10,000$ pixel values (for a $100 \times 100$ pixel image). If the images are faces, the principal components are called eigen-heads because of the relationship between principal component analysis and eigenvalues.

Principal component analysis refers to finding principal components and using them for approximation. In numerical analysis and some parts of statistics, this is called the singular value decomposition, abbreviated $S V D$. This is also "related to" (the same as, once you figure out the definitions and notation) as the Karhunen Loeve expansion.

## Matrix formulation, linear algebra

The time series numbers $X_{t, k}$ for $1 \leq t \leq T$ form the components of a column vector with $T$ components

$$
X_{k}=\left(\begin{array}{c}
X_{1, k} \\
\vdots \\
X_{T, k}
\end{array}\right)
$$

These vectors form the columns of a $T \times n$ data matrix $X$, which may be written

$$
X=\left(\begin{array}{ccccc}
\mid & & \mid & & \mid \\
X_{1} & \cdots & X_{k} & \cdots & X_{n} \\
\mid & & \mid & & \mid
\end{array}\right)
$$

Suppose $U_{t, j}$ are the numbers in principal component $j$. These are the components of $U_{j}$, which is a column vector with $T$ entries

$$
U_{j}=\left(\begin{array}{c}
U_{1, j} \\
\vdots \\
U_{T, j}
\end{array}\right)
$$

Similarly, the residual vector for time series $k$ is

$$
R_{k}=\left(\begin{array}{c}
R_{1, k} \\
\vdots \\
R_{T, k}
\end{array}\right)
$$

The residual vectors form the columns of the residual matrix, as with the data vectors and data matrix:

$$
R=\left(\begin{array}{ccccc}
\mid & & \mid & & \mid \\
R_{1} & \cdots & R_{k} & \cdots & R_{n} \\
\mid & & \mid & & \mid
\end{array}\right)
$$

The expression residuals (1) may be written in vector form as (with $U=U_{1}$ being just the first principal component)

$$
\begin{equation*}
R_{k}=X_{k}-w_{k} U_{1} \tag{3}
\end{equation*}
$$

This may be written in matrix form using a $n$ component row vector

$$
w=\left(\begin{array}{lll}
w_{1} & \cdots & w_{n}
\end{array}\right) .
$$

The matrix form of (11) and (3) is

$$
\begin{equation*}
R=X-U_{1} w \tag{4}
\end{equation*}
$$

This may be rewritten as $X=U_{1} w+R$. In matrix form, this represents the tallthink data matrix $X$ as the sum of the explanation part, which is the column vector $U_{1}$ multiplied by the row vector $w$ and the residual $R$, which has the
same shape as $X$ :


The $(t, k)$ entry of $U_{1} w$ is $U_{t, 1} w_{k}$, which is the prediction of $X_{t, k}$ in (1).
If we use two principal components to "explain" the data matrix $\bar{X}$, we can put the two principal component vectors into a $T \times 2$ matrix $U$. We can put the two sets of weight vectors (row vectors) into a $2 \times n$ matrix $w$. The $T \times 2$ matrix $U$ can multiply the $2 \times n$ matrix $w$, and gives a $T \times n$ matrix $U w$. The residual is $X-U w$. The $(t, k)$ component of the residual is

$$
R_{t, k}=X_{t, k}-\left(U_{t, 1} w_{1, k}+U_{t, 2} w_{2, k}\right)
$$

The matrix form, with shapes, is:


If you use $m$ principal components, the task is to find a $T \times m$ matrix $U$ and an $m \times n$ matrix $w$ so that $R=X-U w$ is minimized, in the least squares sense (2).

A more abstract point of view uses the linear concept of rank of a matrix. If $Y$ is a $T \times n$ matrix, then $\operatorname{rank}(Y)$ is the dimension of the vector space spanned by the columns of $Y$. If $Y$ is tall and thin (as a data matrix may be), then the rank of $Y$ is $n$ if the column vectors are linearly independent. The data matrix $X$ probably usually has this property. But the $T \times n$ matrix $U_{1} w$ has rank 1. To see this, note that the first column of $U_{1} w$ is $U_{1} w_{1}$ and the second column is $U_{1} w_{2}$, etc. Each column is a multiple of the first column (assuming $\left.w_{1} \neq 0\right)$. It is "easy to see" that any rank 1 matrix may be written in this way, as the product of a column vector and a row vector. Therefore, the abstract formulation of the first principal component component problem is to find the rank 1 matrix that best approximates $X$ in the sum of squares sense. With $m$ principal components, the problem is to find the rank $m$ matrix that best approximates $X$.

Suppose $U$ is a column vector and $w$ is a row vector and we know $U w$. Then we know $U$ and $w$ only "up to a constant factor". That is, $2 U_{1}$ and $\frac{1}{2} w$ have the same product as $U w$. We cannot tell by looking at $U w$ whether the column vector is $U$ or $2 U$, because the row vector might be $w$ or $\frac{1}{2} w$. The problem is worse for higher rank approximations. If $U$ is a $T \times m$ matrix and $w$ is an $m \times n$ matrix, then the approximation to $X$ is $U w$. If $A$ is any invertible $m \times m$ matrix, then

$$
U w=(U A)\left(A^{-1} w\right)
$$

If $\widetilde{U}=U A$, then the best approximation with $U$ is the same as the best approximation with $\widetilde{U}$. This is because the "column space" spanned by the columns of $U$ is the same as the column space spanned by the columns of $\widetilde{U}=U A$. The singular value decomposition of $X$ determines singular vectors $U_{k}$ in some order so they become (nearly) uniquely defined.

## First singular value and vector

You can approach the singular value decomposition from a point of view that at first seems different from optimal approximation. The two approaches lead to the same vectors and the same approximations. To find the singular value decomposition, we start by asking about the "maximum stretch" that can be achieved by the matrix $X$. The stretch is measured in the sum of squares sense, in which the "length" of a vector (usually called the vector norm), is

$$
\|v\|=\sqrt{\sum_{k=1}^{n} v_{k}^{2}}
$$

The double bar notation on the left $\|\cdot\|$ is supposed to make norm look like absolute value but fancier (hence, two bars instead of one). The absolute value measures the size of a number and the norm measures the size of a vector. There
are other ways to measure the size of a vector (the largest element, in absolute value, the sum of absolute values, etc.), but this is the one that leads to the singular value decomposition.

If $v$ is an $n$ component column vector, then $y=X v$ is a $T$ component column vector. The stretch is $\|y\| /\|v\|$. The maximum stretch is

$$
\begin{equation*}
\sigma_{1}=\max \|y\|, \quad \text { with }\|v\|=1 \tag{5}
\end{equation*}
$$

This is a constrained optimization problem like the ones we did in mean variance analysis. It is simpler without the square roots, so we solve the equivalent problem

$$
\sigma_{1}^{2}=\max \|y\|^{2}, \quad \text { with }\|v\|^{2}=1
$$

It helps to write the norms as inner products, so

$$
\|y\|^{2}=y_{1} \cdot y_{1}+\cdots+y_{n} \cdot y_{n}=\left(\begin{array}{lll}
y_{1} & \cdots & y_{n}
\end{array}\right)\left(\begin{array}{c}
y_{1} \\
\vdots \\
y_{n}
\end{array}\right)=y^{t} y
$$

With this, the optimization problem is

$$
\begin{equation*}
\sigma_{1}^{2}=\max y^{t} y, \quad \text { with } v^{t} v=1 \tag{6}
\end{equation*}
$$

This becomes more explicit when we substitute $y=X v$ (recall that the transpose of a product of matrices is the product of the transposes, with the order reversed, and that matrix multiplication is associative):

$$
y^{t} y=(X v)^{t}(X v)=\left(v^{t} X^{t}\right)(X v)=v^{t}\left(X^{t} X\right) v
$$

Therefore,

$$
y^{t} y=v^{t} C v, \quad \text { where } C=X^{t} X
$$

The matrix $C=X^{t} X$ is $n \times n$, being the product of $n \times T$ and $T \times n$ matrices:


It is symmetric, because (using the rules of transposes, with $\left.\left(X^{t}\right)^{t}=X\right)$

$$
C^{t}=\left(X^{t} X\right)^{t}=X^{t}\left(X^{t}\right)^{t}=X^{t} X=C
$$

The entries of $C$ are (the $(i, t)$ entry of $X^{t}$ is the $(t, i)$ index of $X$ )

$$
\begin{aligned}
C_{i j} & =\sum_{t=1}^{T}\left(X^{t}\right)_{i t} X_{t j} \\
& =\sum_{t=1}^{T} X_{t i} X_{t j}
\end{aligned}
$$

This is the same as the covariance we saw in mean/variance analysis, except that we have not subtracted the means from the data columns. It is clear from this formula that $C_{i j}=C_{j i}$, which is to say that $C$ is symmetric.

The maximization problem for maximum stretch is

$$
\begin{equation*}
\max v^{t} C v, \quad \text { with } v^{t} v=1 \tag{7}
\end{equation*}
$$

This constrained optimization problem can be solved using a Lagrange multiplier, as in mean/variance analysis. The objective function is $f(v)=v^{t} C v$ and the constraint function (there is just one constraint here) is $g(v)=v^{t} v$. The Lagrange multiplier condition is

$$
\nabla f(v)=\lambda \nabla g(v)
$$

We saw, in mean/variance analysis, that

$$
\nabla v^{t} C v=2 C v
$$

If we take $C=I$ (the identity matrix), we get

$$
\nabla v^{t} v=2 v
$$

This may be seen directly. The $j$ component of $\nabla v^{t} v$ is

$$
\frac{\partial}{\partial v_{j}} \sum_{i=1}^{n} v_{i}^{2}=\sum_{i=1}^{n} \frac{\partial v_{i}^{2}}{\partial v_{j}}=2 v_{j}
$$

Thus, the Lagrange multiplier equation for this problem is

$$
\begin{equation*}
C v=\lambda v \tag{8}
\end{equation*}
$$

An equation of the form (8) is an eigenvalue equation. The the vector $v$ is an eigenvector and the number $\lambda$ is an eigenvalue ${ }^{1}$ The eigenvalue equation

[^0]says that the vector $C v$ is in the same direction as $v$. For most matrices and most vectors, $C v$ is in a different direction. Eigenvector directions $v$ are special. There is a lot of theory and many applications of eigenvalues and eigenvectors. I mention only the things that are needed for singular values and principal component analysis.

The eigenvalue equation (8) may be thought of as an algebraic equation for $\lambda$. It turns out that there is a polynomial $p(\lambda)$ so that the zeros of $p$ (the values of $\lambda$ with $p(\lambda)=0$ ) are eigenvalues. Therefore, every matrix has at least one eigenvalue, if you allow complex eigenvalues. But the eigenvalue $\lambda$ in (8) is not complex. We know about it because of the Lagrange multiplier theory. The symmertic matrix $C$ has at least one real eigenvalue and corresponding real eigenvector. These come from the optimization problem (7). (You might ask why this does not apply to show that any matrix has at least one real eigenvalue. The answer is that $\nabla\left(v^{t} C v\right)=2 C v$ only if $C$ is symmetric. Otherwise, $\nabla(C v)=$ $\left(C+C^{t}\right) v$. This is the same if $C^{t}=C$, but not otherwise. The rotation matrix

$$
C=\left(\begin{array}{cc}
\cos (\theta) & \sin (\theta) \\
-\sin (\theta) & \cos (\theta)
\end{array}\right)
$$

rotates a two component vector by angle $\theta$. A rotated vector $v$ points in a different direction, so the eigenvalue equation (8) cannot be satisfied. Note that

$$
C+C^{t}=\left(\begin{array}{cc}
2 \cos (\theta) & 0 \\
0 & 2 \cos (\theta)
\end{array}\right)
$$

This matrix is symmetric, but it isn't a rotation matrix. The eigenvalues are $e^{i \theta}$ and $e^{-i \theta}$. They are not real for most $\theta$ values.)

If $v$ satisfies the Lagrange multiplier/eigenvalue equation (8), the corresponding value is

$$
v^{t} C v=v^{t}(C v)=v^{t} \lambda v=\lambda v^{t} v=\lambda
$$

The last equation comes from the constraint $v^{t} v=1$. The maximum stretch, according to (6) is

$$
\sigma_{1}=\sqrt{\lambda}
$$

The eigenvalue $\lambda$ in (8) is not unique. Any zero of $p(\lambda)$ is an eigenvalue. The singular value $\sigma_{1}$ is the maximum stretch. It corresponds (as we will see) to the largest eigenvalue of $C$. We call the corresponding eigenvector $v_{1}$. This anticipates that there will be more eigenvectors $v_{2}, \ldots$.

## More principal components and singular values

The vector $v_{1}$ is the first right principal component ("right" because $v$ is on the right of $X$ in $X v$ ). The corresponding left principal component is

$$
U_{1}=\sigma_{1} X v_{1}
$$

The factor $\sigma_{1}$ on the left is so that $\left\|U_{1}\right\|=1$. The vector $v_{1}$ has (we saw above) $\left\|X v_{1}\right\|=\sigma_{1}$, so when we take out the factor of $\sigma_{1}$, the vector $U_{1}$ has length 1 .

Orthogonality of vectors is important for understanding the rest of the principal components and singular values. Column vectors $v$ and $w$ are orthogonal if $v^{t} w=0$. Orthogonality comes up when you minimize or maximize lengths of vectors or distances - if those lengths and distances are measured in the RMS (root mean square) sense $\|v\|=\sqrt{\sum v_{j}^{2}}$. The "mean" square (average square) would be $\frac{1}{n} \sum v_{j}^{2}$ and the "root" mean square would be the square root of that. We call it RMS even factor $\frac{1}{n}$ is left out. A simple form of the minimization property comes from this calculation. It is like the binomial calculation $(a+b)^{2}$, except that you have to keep the of order in matrices, except $v^{t} w=w_{t} v$ :

$$
\begin{aligned}
\frac{d}{d t}\|v+t w\|^{2} & =\frac{d}{d t}\left[(v+t w)^{t}(v+t w)\right] \\
& =\frac{d}{d t}\left[v^{t}(v+t w)+t w^{t}(v+t w)\right] \\
& =\frac{d}{d t}\left[v^{t} v+t v^{t} w+t w^{t} v+t^{2} w^{t} w\right] \\
& =\frac{d}{d t}\left[v^{t} v+2 t v^{t} w+t^{2} w^{t} w\right] \\
& =2 v^{t} w+2 t w^{t} w
\end{aligned}
$$

Set $t=0$ and you get

$$
\left.\frac{d}{d t}\|v+t w\|^{2}\right|_{t=0}=2 v^{t} w=0, \quad \text { if } v^{t} w=0
$$

You can interpret this to say that if you start at $v$ and start moving in the $w$ direction, and if $w$ is orthogonal to $v$, then in the first derivative approximation, at the start, you don't change the length. You can see this without calculus using the calculation we just did and $\|v\|^{2}=v^{t} v$, etc.:

$$
\left\|v^{t} w\right\|=\|v\|^{2}+2 t v^{t} w+t^{2}\|w\|^{2}
$$

This is a quadratic function of $t$. If $v^{t} w=0$, the minimum is at $t=0$. Important for what we're about to do: if $v^{t} w \neq 0$, then $t=0$ is not the minimum. This is the principle: optimizing (minimizing or maximizing) distance implies orthogonality.

Now we define the second principal component $v_{2}$, which is the second singular vector. We also find the second singular value $\sigma_{2}$. Recall that $v_{1}$ and $\sigma_{1}$, and the left principal component $U_{1}$ are defined by

$$
\begin{array}{rlrl}
\sigma_{1} & =\max _{v}\|X v\|, \quad \text { with }\|v\|=1 \\
v_{1} & =\arg \max _{v}\|X v\|, \quad \text { with }\|v\|=1 \\
\sigma_{1} U_{1} & =X v_{1} . &
\end{array}
$$

The next singular value and principal components are defined by solving the same maximization problem with the extra constraint that $v$ is orthogonal to $v_{1}$ :

$$
\begin{array}{rlrl}
\sigma_{2} & =\max _{v}\|X v\|, \quad \text { with } v_{1}^{t} v=0, \quad\|v\|=1 \\
v_{2} & =\arg \max _{v}\|X v\|, \quad \text { with } v_{1}^{t} v=0, \quad\|v\|=1 \\
\sigma_{2} U_{2} & =X v_{2} . &
\end{array}
$$

The important orthogonality thing that happens here is that $U_{2}$ comes out to be orthogonal to $U_{1}$. We made the second right principal component orthogonal to the first one ( $v_{2}$ orthogonal to $v_{1}$ ) by definition. It is a theorem (hopefully not obvious at this point) that the left principal components are also orthogonal: $U_{1}^{t} U_{2}=0$.

We prove the $U_{1}, U_{2}$ orthogonality theorem by showing that if $U_{1}^{t} U_{2} \neq 0$ then $v_{1}$ was not optimal. In fact, if $v^{t} v_{1}=0$, and if $v_{1}$ is optimal, then $X v$ is orthogonal to $U_{1}=X v_{1}$. The proof is by contradiction. We define vectors $Y=X v$ and $Y_{1}=X v_{1}$ to be the vectors $U$ and $U_{1}$ without normalization. The condition $Y^{t} Y_{1} \neq 0$ is the same as $U^{t} U_{1} \neq 0$. If there is a $v$ with $v^{t} v_{1}=0$ and $Y^{t} Y_{1} \neq 0$, then $v_{1}$ is not optimal. The proof is a calculation using the "trial vector" $v(t)=v_{1}+t v$. We know that $v(0)=v_{1}$ has $\|v(0)\|=0$. We also know that, in the first derivative approximation, $\|v(t)\| \approx 1$ also (for small $t$ ). Now define $Y(t)=X v(t)=Y_{1}+t Y$. The same calculation shows that, in the first derivative approximation, $\|Y(0)\|=\left\|Y_{1}\right\|$ is not the maximum value of $\|Y(t)\|$. In fact,

$$
\left.\frac{d}{d t}\|Y(t)\|^{2}\right|_{t=0}=2 Y^{t} Y_{1} \neq 0
$$

At least in the first derivative approximation, $v(0)$ does not optimize the length $\|U(t)\|$ with the constraint $\|v(t)\|=1$.

This argument using first derivative approximations may be a little too informal. There is a more formal proof that uses the same ideas. As often happens, the formal proof hides the idea of the proof inside some more complicated calculations. For the rigorous argument, define a family of trial vectors

$$
v(t)=\frac{1}{\sqrt{1+t^{2}\|v\|^{2}}}\left(v_{1}+t v\right)
$$

In the first derivative approximation, this is the same as $v(t)=v_{1}+t v$ because

$$
\left.\frac{d}{d t} \frac{1}{\sqrt{1+t^{2}\|v\|^{2}}}\right|_{t=0}=0
$$

This is the same as saying that in the first derivative approximation, this "normalization factor" is equal to one and may be left out. The more complicated $v(t)$ has the property that it is always "normalized" to satisfy the constraint
$\|v(t)\|=1$. In fact, using calculations we did just before, and $v^{t} v_{1}=0$, we calculate

$$
\|v(t)\|^{2}=\frac{1}{1+t^{2}\|v\|^{2}}\left\|v_{1}+t v\right\|^{2}=\frac{1}{1+t^{2}\|v\|^{2}}\left(1+t^{2}\|v\|^{2}\right)=1
$$

(The normalization factor $1 / \sqrt{1+t^{2}\|v\|^{2}}$ is chosen to make this happen.)
Now, suppose that $Y=X v$ is not orthogonal to $Y_{1}=X v_{1}$, which is $Y^{t} Y_{1} \neq$ 0 , as before. Then define $Y(t)=X v(t)$ and calculate

$$
\begin{aligned}
\|Y(t)\|^{2} & =\left\|\frac{1}{\sqrt{1+t^{2}\|v\|^{2}}}\left(X v_{1}+t X v\right)\right\|^{2} \\
& =\frac{1}{1+t^{2}\|v\|^{2}}\left\|Y_{1}+t Y\right\|^{2} \\
& =\frac{1}{1+t^{2}\|v\|^{2}}\left(\left\|Y_{1}\right\|^{2}+2 t Y^{t} Y_{1}+t^{2}\|Y\|^{2}\right)
\end{aligned}
$$

The value of $\|Y(0)\|^{2}$ is $\sigma_{1}^{2}$ (by definition of $Y_{1}$ and $\sigma_{1}$ ). The derivative (a calculation any math finance student can do) is

$$
\begin{aligned}
\left.\frac{d}{d t}\|Y(t)\|^{2}\right|_{t=0} & =\left.\frac{d}{d t}\left[\frac{1}{1+t^{2}\|v\|^{2}}\left(\left\|Y_{1}\right\|^{2}+2 t Y^{t} Y_{1}+t^{2}\|Y\|^{2}\right)\right]\right|_{t=0} \\
& =2 Y^{t} Y_{1} \\
& \neq 0
\end{aligned}
$$

This calculation is a combination of the first derivative approximations we just did, plus the product rule (which is really a property of first derivative approximations). It's a harder way of doing the less formal reasoning above. Anyway, if the derivative at $t=0$ is not zero, then $t=0$ is not a local maximum. There are nearby $t$ values with $\|Y(t)\|^{2}>\sigma_{1}^{2}$. This is the proof by contradiction: if $Y X v$ is not orthogonal to $Y_{1}=X v_{1}$, then $\sigma_{1}$ is not the maximum stretch and $v_{1}$ does not optimize the stretch.

In this way, we find vectors $v_{2}$ and $U_{2}=\frac{1}{\sigma_{2}} X v_{2}$ so that $\left\|v_{2}\right\|=1,\left\|U_{2}\right\|=$ $1, v_{2}^{t} v_{1}=0$ and $U_{2}^{t} U_{1}=0$. It is possible to continue and find $v_{3}, v_{4}$, etc. We constrain $v_{3}$ to be orthogonal to $v_{1}$ and $v_{2}$. The vector $Y_{3}=X v_{3}$ will automatically be orthogonal to $Y_{1}$ and $Y_{2}$. If $Y_{3}$ is not orthogonal to $Y_{1}$, then $v_{1}$ was not optimal (as we just saw). If $Y_{3}$ is not orthogonal to $Y_{2}$, then $v_{2}$ is not optimal (same reason). In this way, we construct $v_{j}$ with

$$
\begin{equation*}
\sigma_{j} U_{j}=X v_{j} \tag{9}
\end{equation*}
$$

This process has to stop at $v_{n}$. If $v_{1}, \ldots, v_{n}$ are all orthogonal to each other then there are no $n$ component vectors orthogonal to all of them. Thus, there cannot be a $v_{n+1}$. If $T<n$, the process would have to stop earlier, because the
$Y_{k}$ could not all be orthogonal and non-zero. For now, we continue to suppose $T>n$ (the data matrix $X$ is tall and thin).

To summarize the results, there are $n$ (right principal component) vectors $v_{j}$ so that $\left\|v_{j}\right\|=1$ and if $i \neq j$ then $v_{i}^{t} v_{j}=0$. A family of vectors like this is ortho-normal. They are "ortho" because they are orthogonal to each other. They are "normal" because they are normalized to have length 1 in the RMS sense. There are positive principal stretches $\sigma_{1} \geq \sigma_{2} \geq \cdots \geq \sigma_{n}$. There are $n$ (left principal component) vectors $U_{j}$, which are also ortho-normal, so that

$$
\begin{equation*}
X v_{j}=\sigma_{j} U_{j} \tag{10}
\end{equation*}
$$

## Best least squares approximation

Suppose $z$ and $v$ are $n$ component column vectors, and you want to find the best approximation of $z$ in the form $w v$. Here $w$ is a number, which is a scaling factor. The residual for this approximation is $R=z-w v$. The least squares best approximation minimizes the sum of squares of the residuals. This is a one variable $(w)$ minimization problem that can be solved with ordinary calculus. Here is a vector version of the same calculation:

$$
\begin{aligned}
\|R\|^{2} & =\|z-w v\|^{2} \\
& =(z-w v)^{t}(z-w v) \\
& =z^{t} z-2 w z^{t} v+w^{2} v^{t} v \\
& =\|z\|^{2}-2 w z^{t} v+w^{2}\|v\|^{2} .
\end{aligned}
$$

The derivative with respect to $w$ is $-2 z^{t} v+2 w\|v\|^{2}$. Setting this to zero gives the optimal $w$ :

$$
w_{*}=\frac{z^{t} v}{\|v\|^{2}}
$$

If $v$ is normalized to $\|v\|=1$, then this simplifies to

$$
\begin{equation*}
w_{*}=z^{t} v \tag{11}
\end{equation*}
$$

This approximation has a geometrical interpretation The approximation $z \approx$ $w v$ is the projection (or orthogonal projection) of $z$ onto the line generated by $v$. The residual $R$ is the projection of $z$ onto the plane (or hyperplane, for more than 3 dimensions) orthogonal (perpendicular) to $v$. To see that the residual is perpendicular to $v$, calculate (with $\|v\|^{2} \quad v^{t} v=1$, and $\left.z^{t} v=v^{t} w\right)$ )

$$
v^{t} R=v^{t}(z-w v)=v^{t} z-w v^{t} v=v^{t} z-w=0
$$

This is a form of the orthogonality principle presented earlier: minimum error in the least squares sense leads to orthogonality. Here, the residual is orthogonal to the fitting "function" (the vector $v$ ).

The best approximation using two or more vectors has a similar structure. Suppose we want to minimize the residual using two fitting vectors

$$
R=z-w_{1} v_{1}-w_{2} v_{2}
$$

Suppose the fitting vectors $v_{1}$ and $v_{2}$ are normalized and orthogonal to each other: $\left\|v_{1}\right\|=1,\left\|v_{2}\right\|=1$, and $v_{1}^{t} v_{2}=0$. The residual calculation we did for one vector fitting may be repeated in almost the same form:

$$
\begin{aligned}
\|R\|^{2} & =R^{t} R \\
& =\left(z-w_{1} v_{1}-w_{2} v_{2}\right)^{t}\left(z-w_{1} v_{1}-w_{2} v_{2}\right) \\
& =\|z\|^{2}-2 w_{1} z^{t} v_{1}-w_{2} z^{t} v_{2}+w_{1}^{2}+w_{2}^{2} .
\end{aligned}
$$

This depends independently on $w_{1}$ and $w_{2}$. Minimizing over $w_{1}$ and $w_{2}$ gives

$$
w_{* 1}=v_{1}^{t} z, \quad w_{* 2}=v_{2}^{t} z
$$

The residual (by a similar calculation) is orthogonal to $v_{1}$ and $v_{2}$.
It is common to talk about explained and unexplained, or residal sum of squares. The total sum of squares of $z$ is (in different notations)

$$
S S_{\mathrm{tot}}=\sum_{j=1}^{n} z_{j}^{2}=z^{t} z=\|z\|^{2}
$$

The explained sum of squares is the sum of squares in the approximation $w_{1} v_{1}+$ $w_{2} v_{2}$;. This is (using the fact that $v_{1}$ and $v_{2}$ are ortho-normal)

$$
\begin{align*}
S S_{\exp } & =\left\|w_{1} v_{1}+w_{2} v_{2}\right\|^{2} \\
& =\left(w_{1} v_{1}+w_{2} v_{2}\right)^{t}\left(w_{1} v_{1}+w_{2} v_{2}\right) \\
& =w_{1}^{2} v_{1}^{t} v_{1}+2 w_{1} w_{2} v_{1}^{2} v_{2}+w_{2}^{2} v_{2}^{2} \\
S S_{\exp } & =w_{1}^{2}+w_{2}^{2} . \tag{12}
\end{align*}
$$

The residual sum of squares is

$$
\begin{align*}
S S_{\mathrm{res}} & =\|R\|^{2} \\
& =\left\|z-w_{1} v_{1}-w_{2} v_{2}\right\|^{2} \\
& =\left(z-w_{1} v_{1}-w_{2} v_{2}\right)^{t}\left(z-w_{1} v_{1}-w_{2} v_{2}\right) \\
& =z^{t} z-2 w_{1} z^{t} v_{1}-2 w_{2} z^{t} v_{2}+w_{1}^{2}+w_{2}^{2} \quad\left(\operatorname{using} v_{1}^{t} v_{2}=0, v_{1}^{t} v_{1}=1, \text { etc. }\right) \\
& =z^{t} z-\left(w_{1}^{2}+w_{2}^{2}\right) \\
S S_{\mathrm{res}} & =S S_{\mathrm{tot}}-S S_{\mathrm{exp}} \tag{13}
\end{align*}
$$

In other words, the total sum of squares is the explained plus the unexplained sums of squares. This is the pythagorean theorem of data fitting. The pythagorean
theorem of Pythagoras is about triangles in the plane where one side is orthogonal to another side. The pythagorean theorem of data fitting is about orthonormal fitting vectors and the residual being orthogonal to the best fit. Let $v_{1}$, $\ldots, v_{n}$ be the full set of ortho-normal right principal component vectors. Let $z$ be any $n$ component column vector. The "approximation" of $z$ using all $n$ principal component vectors has no residual (more precisely, $R=0$ ). Thus

$$
\begin{equation*}
z=\sum_{j=1}^{n} w_{j} v_{j}, \quad w_{j}=v_{j}^{t} z \tag{14}
\end{equation*}
$$

All the sum of squares is explained, so

$$
\sum_{i=1}^{n} z_{i}^{2}=\sum_{j=1}^{n} w_{j}^{2}
$$

The singular value decomposition, or $S V D$, of the data matrix $X$ is a representation

$$
\begin{equation*}
X=U \Sigma V^{t} \tag{15}
\end{equation*}
$$

Here $U$ is a $T \times n$ matrix (the same shape as $X$ ) whose columns are the left principal component vectors $U_{j}$. Next, $\Sigma$ is a diagonal $n \times n$ matrix with the singular values $\sigma_{j}$ on the diagonal. Finally, $V$ is an $n \times n$ matrix whose columns are the right principal component vectors $v_{j}$. The SVD equation contains all the orthogonality and approximation information that was used to construct the $v_{j}$. We prove the matrix SVD formula (15) by showing that for any $n$ component column vector $z$,

$$
X z=U \Sigma V^{t} z
$$

The trick is to see how $X$ "acts" on the vector $z$ by knowing how $z$ is represented in terms of the $v_{j}(14)$, and how $X$ acts on $v_{j}(9)$. The result is

$$
\begin{aligned}
X z & =X \sum_{j=1}^{n} w_{j} v_{j} \\
& =\sum_{j=1}^{n} w_{j} X v_{j} \\
& =\sum_{j=1}^{n} \sigma_{j} U_{j} w_{j} \\
& =\sum_{j=1}^{n} \sigma_{j} U_{j} v_{j}^{t} z \\
X z & =\left(\sum_{j=1}^{n} \sigma_{j} U_{j} v_{j}^{t}\right) z
\end{aligned}
$$

This shows that the data matrix may be represented of a sum involving principal components

$$
\begin{equation*}
X=\sum_{j=1}^{n} \sigma_{j} U_{j} v_{j}^{t} \tag{16}
\end{equation*}
$$

This is the main formula of principal component analysis. It implies the optimality properties of the principal components $U_{j}$ mentioned at the beginning (see below). We will see that this is equivalent to the matrix form 15 .

The matrix $U$ has columns $U_{j}$. Define

$$
\left(\begin{array}{ccc}
\mid & \mid & \mid \\
\sigma_{1} U_{1} & \sigma_{j} U_{j} & \sigma_{n} U_{n} \\
\mid & \mid & \mid
\end{array}\right)=\left(\begin{array}{ccc}
\mid & \mid & \mid \\
& & \\
U_{1} & U_{j} & U_{n} \\
& & \mid
\end{array}\right)\left(\begin{array}{cccc}
\sigma_{1} & 0 & \cdots & 0 \\
0 & \sigma_{2} & 0 & \\
\vdots & 0 & \ddots & \vdots \\
& & & \mid \\
0 & \cdots & 0 & \sigma_{n}
\end{array}\right)
$$

The first matrix on the right is $U$. The second is a diagonal matrix called $\Sigma$. Define columns vectors $Y_{j}=\sigma_{j} U_{j}$, and the $T \times n$ matrix $Y$ with columns $Y_{j}$. This says that $Y=U \Sigma$. Therefore, the claim that 16 is equal to 15 is the same as

$$
\left(\begin{array}{ccc}
\mid & \mid & \mid \\
Y_{1} & Y_{j} & Y_{n} \\
\mid & \mid & \mid
\end{array}\right)\left(\begin{array}{ccccc}
- & - & v_{1}^{t} & - & - \\
- & - & v_{j}^{t} & - & - \\
- & - & v_{n}^{t} & - & -
\end{array}\right)=\sum_{j=1}^{n} Y_{j} v_{j}^{t}
$$

We can verify this by looking at the elements of the matrices. We choose to verify it by applying it to a column vector $z$. We define "weights" $w_{j}=v_{j}^{t} z$ as before. On the right side, we get

$$
\left(\sum_{j=1}^{n} Y_{j} v_{j}^{t}\right) z=\sum_{j=1}^{n} Y_{j}\left(v_{j}^{t} z\right)=\sum_{j=1}^{n} w_{j} Y_{j}
$$

On the right, we have

$$
\begin{aligned}
\left(\begin{array}{ccc}
\mid & \mid & \mid \\
Y_{1} & Y_{j} & Y_{n} \\
\mid & \mid & \mid
\end{array}\right)\left(\begin{array}{lllll}
- & - & v_{1}^{t} & - & - \\
- & - & v_{j}^{t} & - & - \\
- & - & v_{n}^{t} & - & -
\end{array}\right) z & =\left(\begin{array}{ccc}
\mid & \mid & \mid \\
Y_{1} & Y_{j} & Y_{n} \\
\mid & \mid & \mid
\end{array}\right)\left(\begin{array}{c}
v_{1}^{t} z \\
\vdots \\
v_{j}^{t} z \\
\vdots \\
v_{n}^{t} z
\end{array}\right) \\
& =\left(\begin{array}{ccc}
\mid & \mid & \mid \\
Y_{1} & Y_{j} & Y_{n} \\
\mid & \mid & \mid
\end{array}\right)\left(\begin{array}{c}
w_{1} \\
\vdots \\
w_{j} \\
\vdots \\
w_{n}
\end{array}\right) \\
& =w_{1} Y_{1}+\cdots+w_{j} Y_{j}+\cdots+w_{n} Y_{n}
\end{aligned}
$$

This is $\sum w_{j} Y_{j}$. Thus, the PCA form (16) is equivalent to the SVD from (15).
We are finally ready to go back to the motivating formula (1). The best joint approximation to the $n$ time series using a single principal component is the first principal component $U_{1}$. In matrix form, we use $\sigma_{1} U_{1} v_{1}^{t}$ as an approximation to $X$. The residual is (see the PCA formula 16)

$$
R=X-\sigma_{1} U_{1} v_{1}^{t}=\sum_{j=2}^{n} \sigma_{j} U_{j} v_{j}^{t}
$$

In components, the first principal component approximation (with residual) is

$$
X_{t, k}=U_{t, 1} \sigma_{1} v_{1, k}+R_{t, k}
$$

Therefore, the approximation weights are $w_{k}=\sigma_{1} v_{1, k}$. The best approximation (in the sum of squares sense) using $m$ time series is

$$
X_{t, k} \approx \sum_{j=1}^{m} w_{j k} U_{t, j}, \quad w_{j k}=\sigma_{j} v_{j, k}
$$

In matrix form, the residual from this approximation is

$$
R=X-\sum_{j=1}^{m} \sigma_{j} U_{j} v_{j}^{t}=\sum_{j=m+1}^{n} \sigma_{j} U_{j} v_{j}^{t}
$$

The total sum of squares for all the series is

$$
S S_{\mathrm{tot}}=\sum_{t=1}^{T} \sum_{j=1}^{n} X_{t j}^{2}=\sum_{j=1}^{n} \sigma_{j}^{2}
$$

Some calculations show that the explained part, using $m$ principal components, is

$$
S S_{\exp }=\sum_{j=1}^{m} \sigma_{j}^{2}
$$

The unexplained, or residual sum of squares is

$$
S S_{\mathrm{res}}=\sum_{t=1}^{T} \sum_{j=1}^{n} R_{t j}^{2}=\sum_{j=m+1}^{n} \sigma_{j}^{2} .
$$


[^0]:    1 The these terms are half translations from the original German, where "eigen" means "proper" and "Werte" means value. The half-translation of normal German "eigenwerte" becomes the weird "eigenvalue". The French did a more complete translation to "valeur propre" (two French words).

