

L9: Principal components analysis

The curse of dimensionality

Dimensionality reduction

Feature selection vs. feature extraction

Signal representation vs. signal classification

Principal components analysis

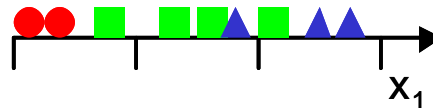
The curse of dimensionality

The curse of dimensionality

- A term coined by Bellman in 1961
- Refers to the problems associated with multivariate data analysis as the dimensionality increases

Consider a 3-class pattern recognition problem

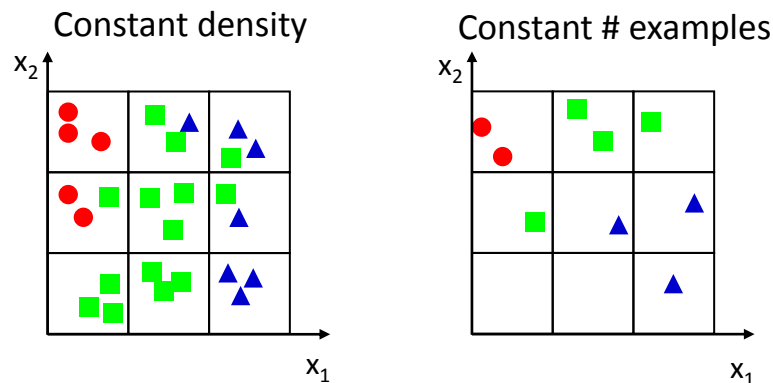
- A simple approach would be to
 - Divide the feature space into uniform bins
 - Compute the ratio of examples for each class at each bin and,
 - For a new example, choose the predominant class in its bin
- In our toy problem we decide to start with one single feature and divide the real line into 3 segments



- After doing this, we notice that there exists too much overlap among the classes, so we decide to incorporate a second feature to try and improve separability

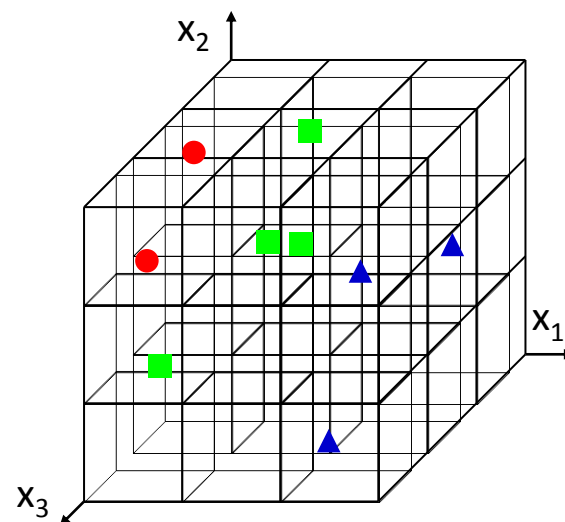
We decide to preserve the granularity of each axis, which raises the number of bins from 3 (in 1D) to $3^2 = 9$ (in 2D)

- At this point we need to make a decision: do we maintain the density of examples per bin or do we keep the number of examples had for the one-dimensional case?
- Choosing to maintain the density increases the number of examples from 9 (in 1D) to 27 (in 2D)
- Choosing to maintain the number of examples results in a 2D scatter plot that is very sparse



Moving to three features makes the problem worse

- The number of bins grows to $3^3 = 27$
- For the same density of examples the number of needed examples becomes 81
- For the same number of examples, the 3D scatter plot is almost empty



Obviously, our approach to divide the sample space into equally spaced bins was quite inefficient

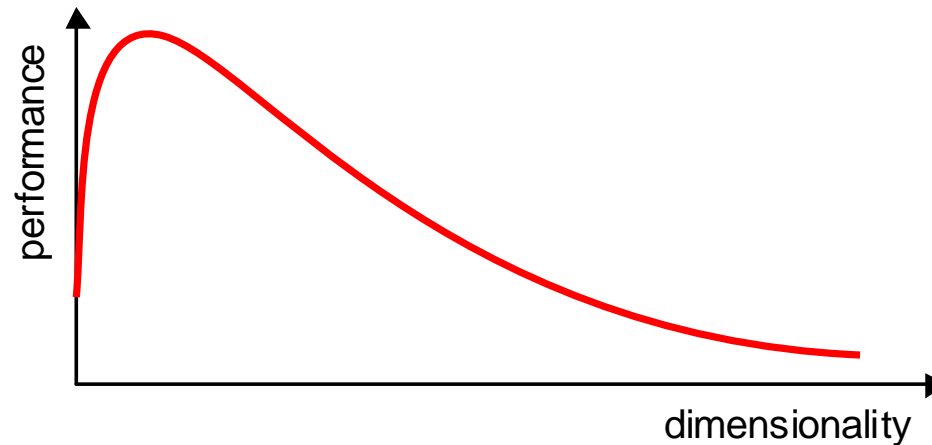
- There are other approaches that are much less susceptible to the curse of dimensionality, **but the problem still exists**

How do we beat the curse of dimensionality?

- By incorporating prior knowledge
- By providing increasing smoothness of the target function
- By reducing the dimensionality

What does the curse of dimensionality mean, in practice?

- For a given sample size, there is a maximum number of features above which the performance of our classifier will degrade rather than improve
- In most cases, the additional information that is lost by discarding some features is (more than) compensated by a more accurate mapping in the lower-dimensional space



Additional implications

- Exponential growth in the #examples required to maintain a given sampling density
 - For a density of N examples/bin and D dimensions, the total number of examples is N^D
- Exponential growth in the complexity of the target function (a density estimate) with increasing dimensionality
 - *“A function defined in high-dimensional space is likely to be much more complex than a function defined in a lower-dimensional space, and those complications are harder to discern”* –J. Friedman
 - This means that, in order to learn it well, a more complex target function requires denser sample points!
- What to do if it ain't Gaussian?
 - For 1D a large number of density functions can be found in textbooks, but for high-dimensions only the multivariate Gaussian density is available.
 - Moreover, for large D the Gaussian can only be handled in a simplified form!
- Humans have an extraordinary capacity to discern patterns and clusters in 1D, 2D and 3D, but these capabilities break down for $D \geq 4$

Dimensionality reduction

Two approaches are available to reduce dimensionality

- **Feature extraction:** creating a subset of new features by combinations of the existing features
- **Feature selection:** choosing a subset of all the features

$$\begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_N \end{bmatrix} \rightarrow \begin{bmatrix} x_{i_1} \\ x_{i_2} \\ \vdots \\ x_{i_M} \end{bmatrix} \quad \begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_N \end{bmatrix} \rightarrow \begin{bmatrix} y_1 \\ y_2 \\ \vdots \\ y_M \end{bmatrix} = f \left(\begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_N \end{bmatrix} \right)$$

The problem of feature extraction can be stated as

- Given a feature space $x_i \in \mathfrak{R}^N$ find a mapping $y = f(x): R^N \rightarrow R^M$ with $M < N$ such that the transformed feature vector $y \in R^M$ preserves (most of) the information or structure in R^N
- An optimal mapping $y = f(x)$ is one that does not increase $P[\text{error}]$
- This is, a Bayes decision rule applied to the initial space R^N and to the reduced space R^M yield the same classification rate

Linear dimensionality reduction

- In general, the optimal mapping $y = f(x)$ will be a non-linear function
 - However, there is no systematic way to generate non-linear transforms
 - The selection of a particular subset of transforms is problem dependent
- For these reasons, feature extraction is commonly based on linear transforms, of the form $y = Wx$

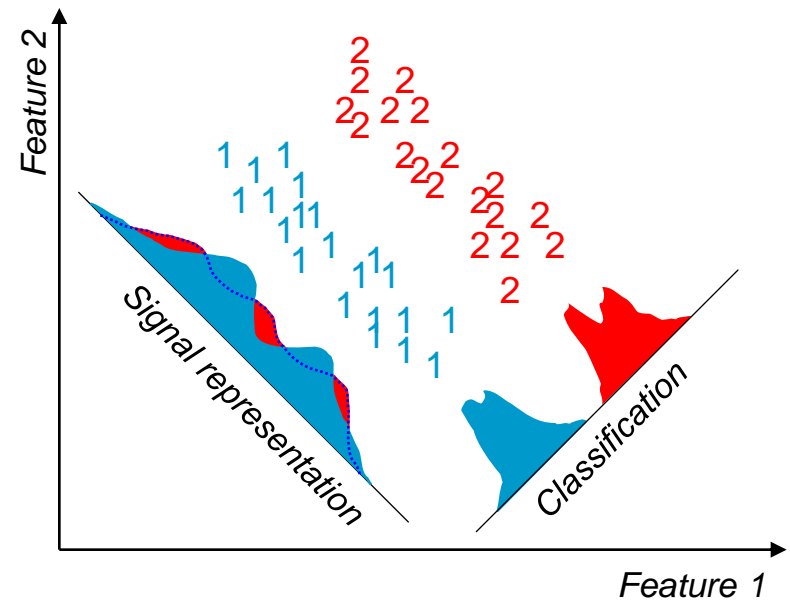
$$\begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_N \end{bmatrix} \rightarrow \begin{bmatrix} y_1 \\ y_2 \\ \vdots \\ y_M \end{bmatrix} = \begin{bmatrix} w_{11} & w_{12} & \dots & w_{1N} \\ \vdots & \vdots & \ddots & \vdots \\ w_{M1} & & & w_{MN} \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_N \end{bmatrix}$$

- NOTE: When the mapping is a non-linear function, the reduced space is called a manifold
- We will focus on linear feature extraction for now, and revisit non-linear techniques when we cover multi-layer perceptrons, manifold learning, and kernel methods

Signal representation versus classification

Finding the mapping $y = f(x)$ is guided by an objective function that we seek to maximize (or minimize)

- Depending on the criteria used by the objective function, feature extraction techniques are grouped into two categories:
 - **Signal representation:** The goal of the feature extraction mapping is to represent the samples accurately in a lower-dimensional space
 - **Classification:** The goal of the feature extraction mapping is to enhance the class-discriminatory information in the lower-dimensional space
- Within the realm of linear feature extraction, two techniques are commonly used
 - Principal components analysis (**PCA**): uses a signal representation criterion
 - Linear discriminant analysis (**LDA**): uses a signal classification criterion



Principal components analysis (PCA)

PCA seeks preserve as much of the randomness (variance) in the high-dimensional space as possible

- Let $x \in \mathfrak{R}^N$ be represented as a linear combination of orthonormal basis vectors $[\varphi_1 | \varphi_2 | \dots | \varphi_N]$ as

$$x = \sum_{i=1}^N y_i \varphi_i \text{ where } \varphi_i^T \varphi_j = \begin{cases} 0; & i \neq j \\ 1; & i = j \end{cases}$$

- Suppose we want to represent x with only M ($M < N$) basis vectors
- We can do this by replacing the components $[y_{M+1}, \dots, y_N]^T$ with some pre-selected constants b_i

$$\hat{x}(M) = \sum_{i=1}^M y_i \varphi_i + \sum_{i=M+1}^N b_i \varphi_i$$

- The representation error is then

$$\begin{aligned} \Delta x(M) &= x - \hat{x}(M) = \\ &= \sum_{i=1}^N y_i \varphi_i - \left(\sum_{i=1}^M y_i \varphi_i + \sum_{i=M+1}^N b_i \varphi_i \right) = \\ &= \sum_{i=M+1}^N (y_i - b_i) \varphi_i \end{aligned}$$

- We can measure this representation error by the mean-squared magnitude of Δx
- Our goal is to find the basis vectors φ_i and constants b_i that minimize this mean-square error

$$\begin{aligned}\bar{\epsilon}^2(M) &= E[|\Delta x(M)|^2] = \\ &E\left[\sum_{i=M+1}^N \sum_{j=M+1}^N (y_i - b_i)(y_j - b_j) \varphi_i^T \varphi_j\right] = \\ &\sum_{i=M+1}^N E[(y_i - b_i)^2]\end{aligned}$$

- To find the optimal values of b_i we compute the partial derivative of the objective function and equate it to zero

$$\begin{aligned}\frac{\partial}{\partial b_i} E[(y_i - b_i)^2] &= -2(E[y_i] - b_i) = 0 \\ &\Rightarrow b_i = E[y_i]\end{aligned}$$

- Therefore, we will replace the discarded dimensions by their expected value, which is an intuitive result

- The MSE can then be written as

$$\begin{aligned}
 \bar{\epsilon}^2(M) &= \sum_{i=M+1}^N E[(y_i - E[y_i])^2] = \\
 &= \sum_{i=M+1}^N E[(x\varphi_i - E[x\varphi_i])^T (x\varphi_i - E[x\varphi_i])] = \\
 &= \sum_{i=M+1}^N \varphi_i^T E[(x - E[x])(x - E[x])^T] \varphi_i = \\
 &= \sum_{i=M+1}^N \varphi_i^T \Sigma_x \varphi_i
 \end{aligned}$$

- We seek the solution that minimizes this expression, subject to the orthonormality constraint, which we incorporate into the expression using a set of Lagrange multipliers λ_i

$$\bar{\epsilon}^2(M) = \sum_{i=M+1}^N \varphi_i^T \Sigma_x \varphi_i + \sum_{i=M+1}^N \lambda_i (1 - \varphi_i^T \varphi_i)$$

- Computing the partial derivative with respect to the basis vectors

$$\begin{aligned}
 \frac{\partial \bar{\epsilon}^2(M)}{\partial \varphi_i} &= \frac{\partial}{\partial \varphi_i} \left[\sum_{i=M+1}^N \varphi_i^T \Sigma_x \varphi_i + \sum_{i=M+1}^N \lambda_i (1 - \varphi_i^T \varphi_i) \right] = \\
 &= 2(\Sigma_x \varphi_i - \lambda_i \varphi_i) = 0 \Rightarrow \Sigma_x \varphi_i = \lambda_i \varphi_i
 \end{aligned}$$

- NOTE: $\frac{d}{dx} (x^T A x) = (A + A^T)x = 2Ax$ (for A symmetric)

- **So φ_i and λ_i are the eigenvectors and eigenvalues of Σ_x**

- We can then express the sum-squared error as

$$\bar{\epsilon}^2(M) = \sum_{i=M+1}^N \varphi_i^T \Sigma_x \varphi_i = \sum_{i=M+1}^N \varphi_i^T \lambda_i \varphi_i = \sum_{i=M+1}^N \lambda_i$$

- In order to minimize this measure, λ_i will have to be smallest eigenvalues
- Therefore, to represent x with minimum MSE, we will choose the eigenvectors φ_i corresponding to the largest eigenvalues λ_i

PCA dimensionality reduction

The optimal* approximation of a random vector $x \in \mathfrak{R}^N$ by a linear combination of $M < N$ independent vectors is obtained by projecting x onto the eigenvectors φ_i corresponding to the largest eigenvalues λ_i of the covariance matrix Σ_x

*optimality is defined as the minimum of the sum-square magnitude of the approximation error

NOTES

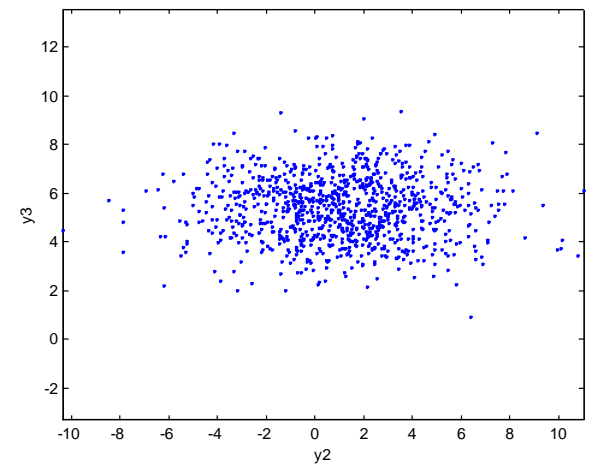
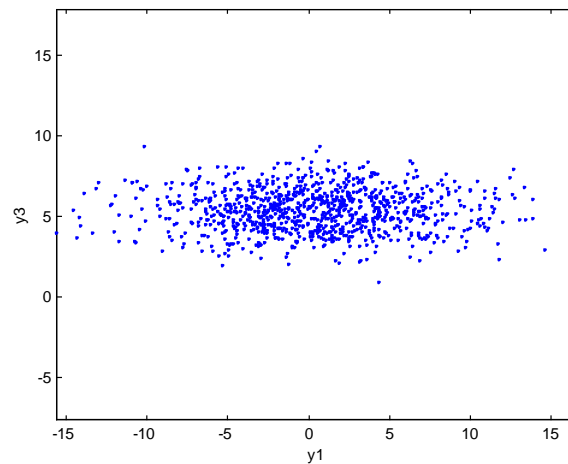
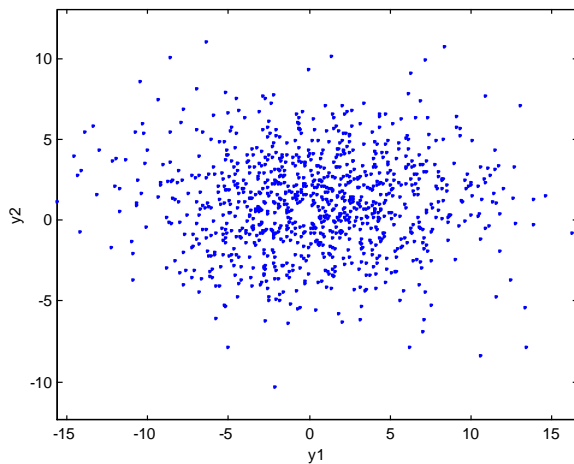
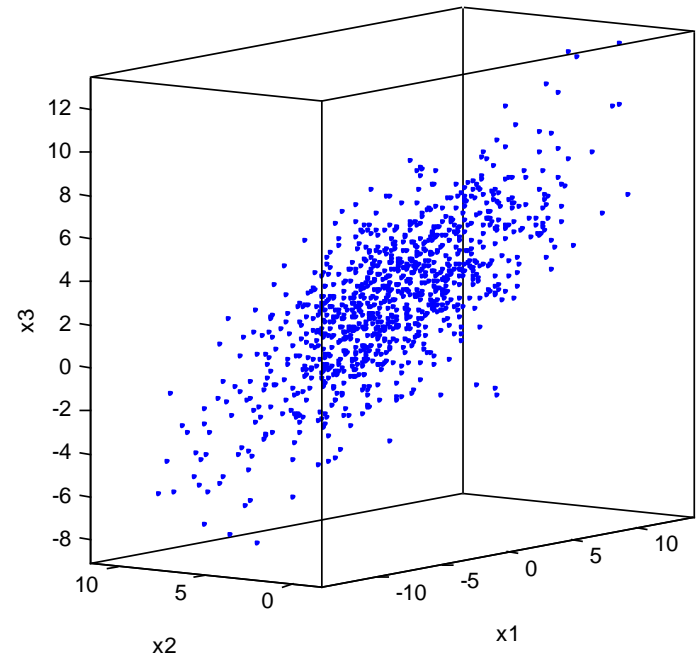
- Since PCA uses the eigenvectors of the covariance matrix Σ_x , it is able to find the independent axes of the data under the unimodal Gaussian assumption
 - For non-Gaussian or multi-modal Gaussian data, PCA simply de-correlates the axes
- The main limitation of PCA is that it does not consider class separability since it does not take into account the class label of the feature vector
 - PCA simply performs a coordinate rotation that aligns the transformed axes with the directions of maximum variance
 - **There is no guarantee that the directions of maximum variance will contain good features for discrimination**
- Historical remarks
 - Principal Components Analysis is the oldest technique in multivariate analysis
 - PCA is also known as the Karhunen-Loève transform (communication theory)
 - PCA was first introduced by Pearson in 1901, and it experienced several modifications until it was generalized by Loève in 1963

Example I

3D Gaussian distribution $N(\mu, \Sigma)$

$$\mu = [0 \ 5 \ 2]^T \text{ and } \Sigma = \begin{bmatrix} 25 & -1 & 7 \\ & 4 & -4 \\ & & 10 \end{bmatrix}$$

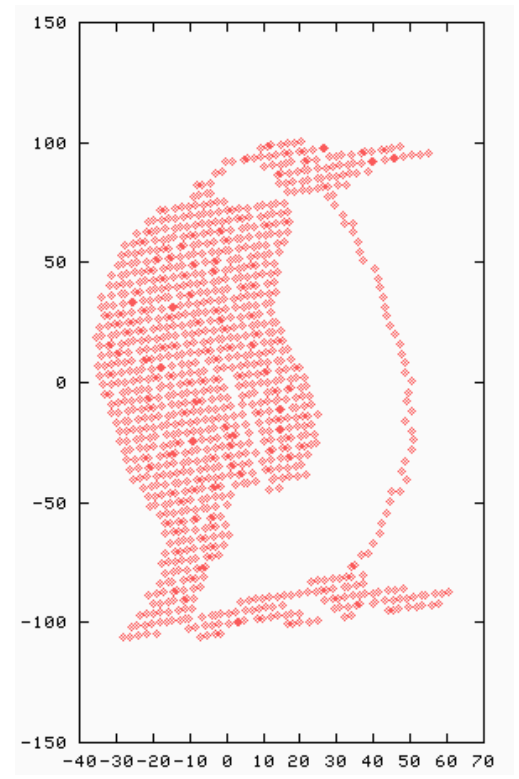
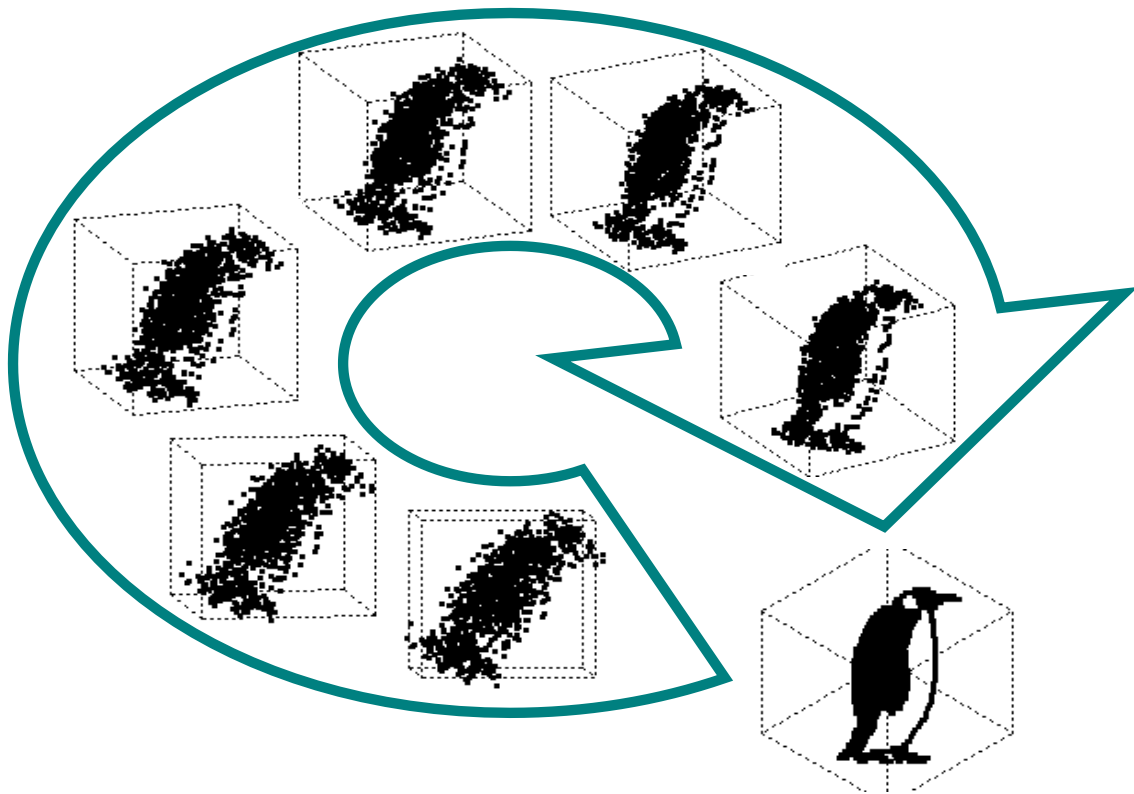
- The three pairs of PCA projections are shown below
 - Notice that PC1 has the largest variance, followed by PC2
 - Also notice how PCA decorrelates the axes



Example II

The example in the next slide shows a projection of a 3D data set into two dimensions

- Initially, there is no apparent structure in the dataset, except for the elongation of the point cloud
- Choosing an appropriate rotation unveils the underlying structure
 - You can think of this rotation as "walking around" the 3D dataset, looking for the best viewpoint
- PCA can help find such underlying structure
 - It selects a rotation such that most of the variability within the data set is represented in the first few dimensions of the rotated data
 - In our 3D case, this may seem of little use
 - However, when the data is highly multidimensional (10's of dimensions), this analysis is quite powerful



Example III

Problem statement

- Compute the PCA for dataset
 $X = \{(1,2), (3,3), (3,5), (5,4), (5,6), (6,5), (8,7), (9,8)\}$
- Let's first plot the data to get an idea of which solution we should expect

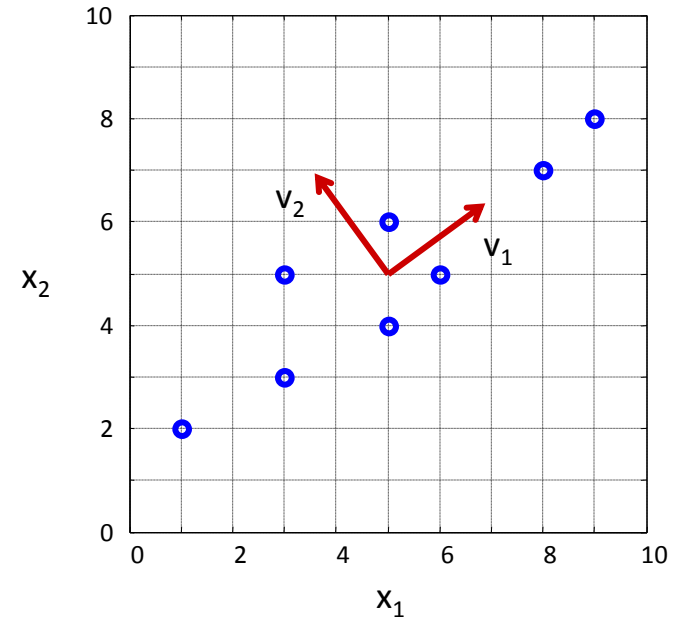
Solution

- The sample covariance is

$$\Sigma_x = \begin{bmatrix} 6.25 & 4.25 \\ 4.25 & 3.5 \end{bmatrix}$$

- The eigenvalues are the zeros of the characteristic equation

$$\begin{aligned} \Sigma_x v &= \lambda v \Rightarrow |\Sigma_x - \lambda I| = 0 \\ \Rightarrow \begin{vmatrix} 6.25 - \lambda & 4.25 \\ 4.25 & 3.5 - \lambda \end{vmatrix} &= 0 \\ \Rightarrow \lambda_1 &= \mathbf{9.34}; \lambda_2 = \mathbf{0.41} \end{aligned}$$



- The eigenvectors are the solutions of the system

$$\begin{bmatrix} 6.25 & 4.25 \\ 4.25 & 3.5 \end{bmatrix} \begin{bmatrix} v_{11} \\ v_{12} \end{bmatrix} = \begin{bmatrix} \lambda_1 v_{11} \\ \lambda_1 v_{12} \end{bmatrix} \Rightarrow \begin{bmatrix} v_{11} \\ v_{12} \end{bmatrix} = \begin{bmatrix} 0.81 \\ 0.59 \end{bmatrix}$$

$$\begin{bmatrix} 6.25 & 4.25 \\ 4.25 & 3.5 \end{bmatrix} \begin{bmatrix} v_{21} \\ v_{22} \end{bmatrix} = \begin{bmatrix} \lambda_2 v_{21} \\ \lambda_2 v_{22} \end{bmatrix} \Rightarrow \begin{bmatrix} v_{21} \\ v_{22} \end{bmatrix} = \begin{bmatrix} -0.59 \\ 0.81 \end{bmatrix}$$

- HINT: To solve each system manually, first assume that one of the variables is equal to one (i.e. $v_{i1} = 1$), then find the other one and finally normalize the vector to make it unit-length