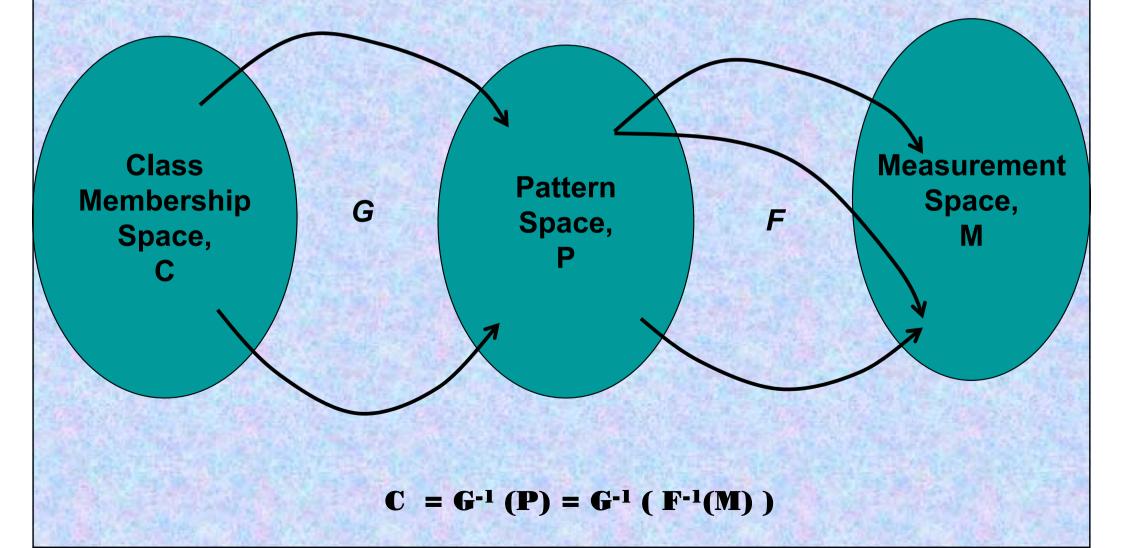
Pattern Recognition

Pattern Recognition is a branch of science that concerns the description or classification (or identification) of measurements. It is an important component of intelligent systems and are used for both data processing and decision making.



Significant Impact •Mathematics

- •Engineering
- •Medicine
- Remote sensing
- •Computer Science

Popular Features Used for analysis

- •Moments
- •Euler Number
- •Chain Code
- Polygonal Approximation
- •Distance Transform
- •B-Spline
- Spectral domain

Application

- Automatic Target Recognition
- Generic Object
 Recognition
- Scene Correlation and Matching
- Landmark Identification from remote sensed data
 - Biometry

An Important Area of Research in Computer Vision and Visual Perception

Features must be invariant to:

- Translation
- Rotation
- Scale
- Noise
- Projective (?)

Computational cost must not be high

Must be distinct and unique for a given shape.

Preferably have graceful degradation due to discontinuities and missing parts

Statistical Features

The features used in pattern recognition and segmentation are generally geometric or intensity gradient based.

One approach is to work directly with regions of pixels in the image, and to describe them by various statistical measures. Such measures are usually represented by a single value. These can be calculated as a simple by-product of the segmentation procedures previously described.

Such *statistical descriptions* may be divided into two distinct classes. Examples of each class are given below:

• Geometric descriptions: area, length, perimeter, elongation, average radius, compactness and moment of inertia.

Topological descriptions: connectivity and Euler number.

Elongation

- sometimes called *eccentricity*. This is the ratio of the maximum length of line or *chord* that spans the region to the minimum length chord. We can also define this in terms of moments, as we will see shortly.

Compactness

- this is the ratio of the square of the perimeter to the area of the region

Connectivity -

- the number of neighboring features adjoining the region.

Euler Number

- for a single region, one minus the number of holes in that region. The Euler number for a set of connected regions can be calculated as the number of regions minus the number of holes.

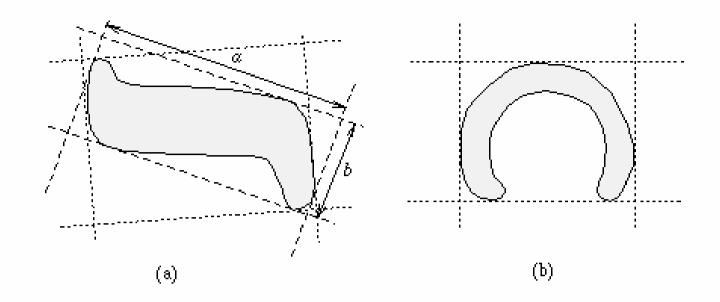


Figure 6.24 Elongatedness: (a) Bounding rectangle gives acceptable results, (b) bounding rectangle cannot represent elongatedness.

Elongatedness:

A ratio between the length and width of the region bounding rectangle = a/b = Area/sqr(thickness).

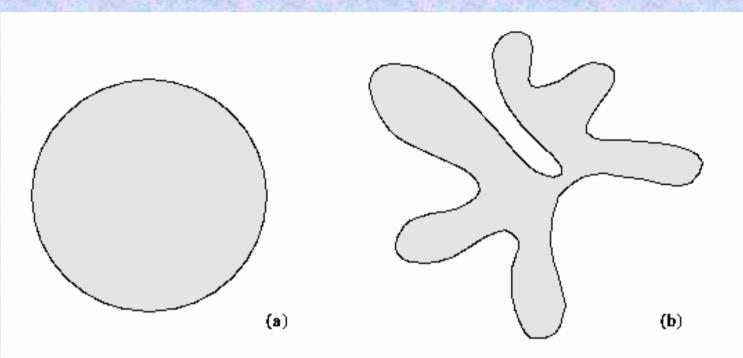


Figure 6.25 Compactness: (a) Compact, (b) non-compact.

Compactness

Compactness is independent of linear transformations

= sqr(perimeter)/Area

Moments of Inertia

The *ij*-th discrete central moment m_{ii}, of a region is defined by:

$$m_{ij} = \sum (x - x)^i (y - y)^j$$

where the sums are taken over all points (x, y) contained within the region and $(x^{,}, y^{-})$ are the center of gravity of the region:

$$\tilde{x} = \frac{1}{n} \sum_{i} x_{i}$$
 and $\tilde{y} = \frac{1}{n} \sum_{i} y_{i}$

Note that, *n*, the total number of points contained in the region, is a measure of its area.

We can form seven new moments from the central moments that are invariant to changes of position, scale and orientation (RTS) of the object represented by the region, although these new moments are *not* invariant under perspective projection.

For moments of order up to seven, these are:

$$\begin{split} M_1 &= m_{20} + m_{02} \\ M_2 &= (m_{20} - m_{02})^2 + 4m_{11}^2 \\ M_3 &= (m_{30} - 3m_{12})^2 + (3m_{21} - m_{03})^2 \\ M_4 &= (m_{30} + m_{12})^2 + (m_{21} + m_{03})^2 \\ M_5 &= (m_{30} - 3m_{12})(m_{30} + m_{12}) \left[(m_{30} + m_{12})^2 - 3(m_{21} + m_{03})^2 \right] \\ &+ (3m_{21} - m_{03})(m_{21} + m_{03}) \left[3(m_{30} + m_{12})^2 - (m_{21} + m_{03})^2 \right] \\ M_6 &= (m_{20} + m_{02}) \left[(m_{30} + m_{12})^2 - 3(m_{21} + m_{03})^2 \right] \\ &+ 4m_{11}(m_{30} + m_{12})(m_{03} + m_{21}) \\ M_7 &= (3m_{21} - m_{03})(m_{12} + m_{30}) \left[(m_{30} + m_{12}^2 - 3(m_{21} + m_{03})^2 \right] \\ &- (m_{30} - 3m_{12})(m_{12} + m_{03}) \left[3(m_{30} + m_{12})^2 - (m_{21} + m_{03})^2 \right] \end{split}$$

We can also define eccentricity, using moments as

$$ext{cccentricity} = rac{m_{20}+m_{02}+\sqrt{(m_{20}-m_{02})^2+4m_{11}^2}}{m_{20}+m_{02}-\sqrt{(m_{20}-m_{02})^2+4m_{11}^2}}.$$

We can also find *principal axes of inertia* that define a natural coordinate system for a region. It is given by:

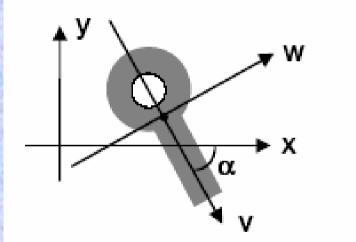
$$heta = rac{1}{2} an^{-1} \left| rac{2m_{11}}{m_{20} - m_{02}}
ight|$$

Geometric properties in terms of moments:

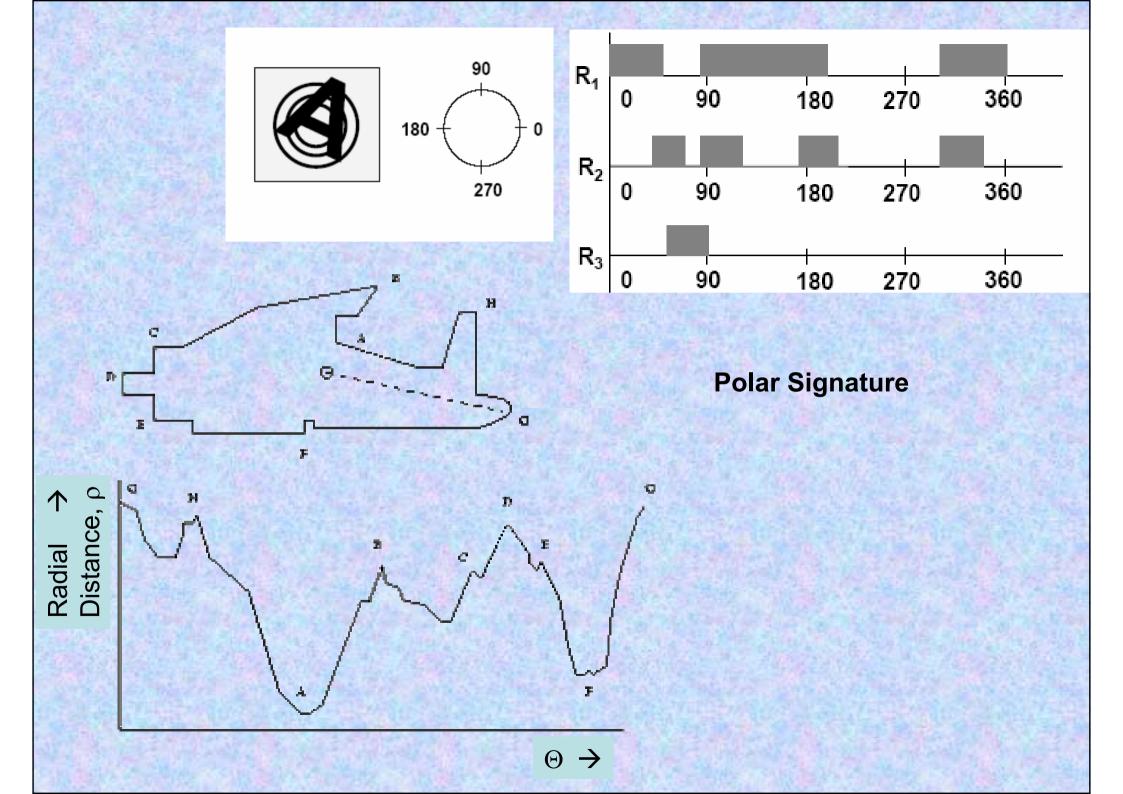
Area =
$$m_{00}$$
; $x = \frac{m_{10}}{m_{00}}$; $y = \frac{m_{01}}{m_{00}}$

Axis of minimal inertia

 $\tan 2\alpha = 2m_{xy}/(m_y - m_x)$



Polar Signature, Skeletons (MAT), B-splines are also used.



Phase of DFT-based Signature Function

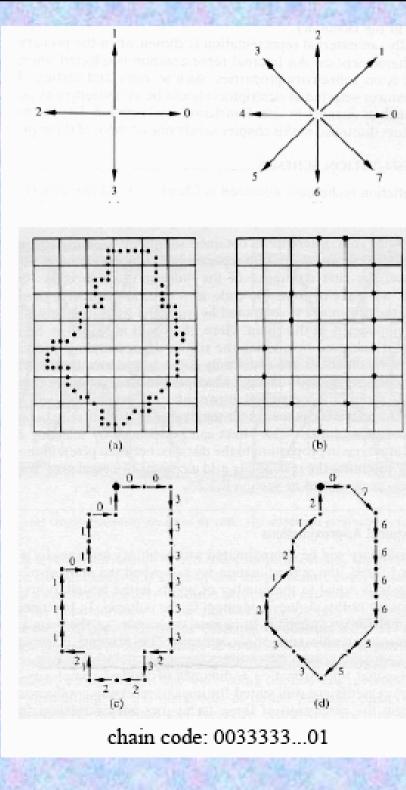
Coordinates of the boundary points of the shape are expressed as: (x_0, y_0) , (x_1, y_1) (x_N, y_N) CR $z_i = x_i + jy_i$

Coefficient of the Fourier Transform is given as $Z(e^{j\omega}) = \sum z_k e^{(-j\omega n)}$ Shape is rotated by an angle θ , and the starting point by I_0 $Z'_m = Z_m e^{j\theta} e^{(-jm2\pi l_0)}$ $Z_m = R_m e^{j\theta_m}$ Shift in phase, due to rotation of the shape and change in starting point is $\theta'_m = \theta_m + \theta - 2\pi l_0 m / N$

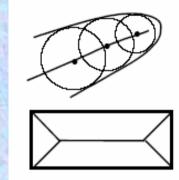
This leads to the derivation of the modified DFT coefficients [1] for normalization against the scaling, rotation and shift in the starting point as depicted in the table below:



Invariance	Modified coefficients
Translation	$Z'_{0} = 0$
Scale	$\mathbf{R'}_{\mathbf{m}} = \mathbf{R}_{\mathbf{m}} / \mathbf{S}$
Rotation	$\Theta'_{m} = \Theta_{m} + \theta - (\Theta_{-1} + \Theta_{+1})/2$
Starting point	$\Theta'_{\rm m} = \Theta_{\rm m} + {\rm m}(\Theta_{-1} - \Theta_{+1})/2$







Skeletons - MAT

Read about :

CSS, a multi-scale representation;

• MCC

- Wavelet based descriptors
- Distance functions Hausdorff
- Shape Context

An Example of Classification

 "Sorting incoming Fish on a conveyor according to species using optical sensing

Sea bass

Species

Salmon

 Some properties that could be possibly used to distinguish between the two types of fishes is

- Length
- Lightness
- Width
- Number and shape of fins
- Position of the mouth, etc...

Features

– This is the set of all suggested features to explore for use in our classifier!

Feature is a property (or characteristics) of an object (quantifiable or non quantifiable) which is used to distinguish between (or classify) two objects.

Feature vector

- A Single feature may not be useful always for classification
- A set of features used for classification form a feature vector

Fish
$$\longrightarrow x^T = [x_1, x_2]$$

Lightness Width

Feature space

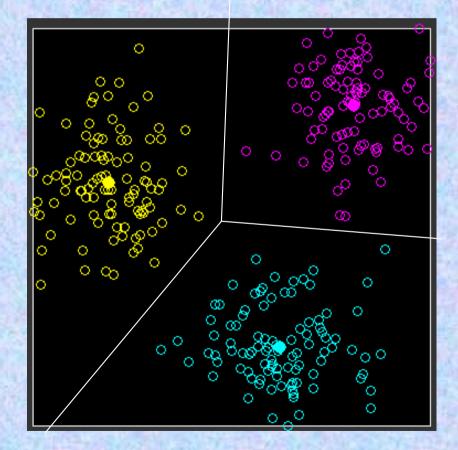
- The samples of input (when represented by their features) are represented as points in the feature space
- If a single feature is used, then work on a one- dimensional feature space.



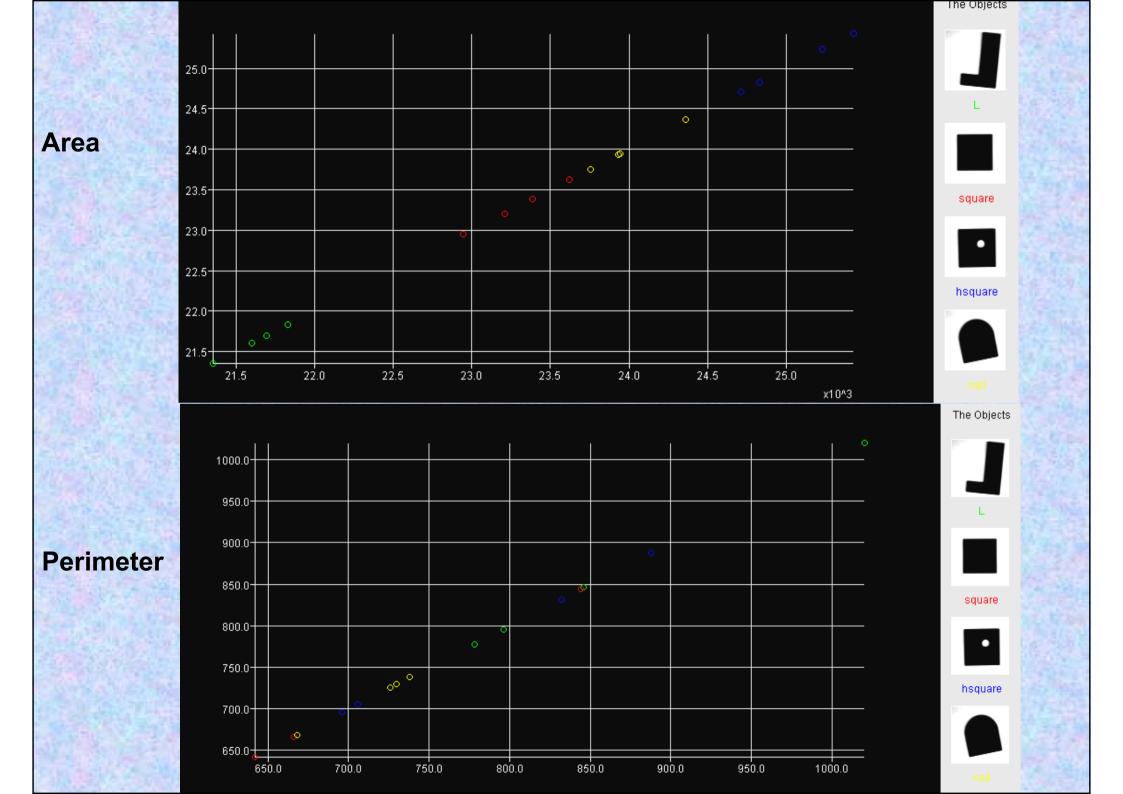
Point representing samples

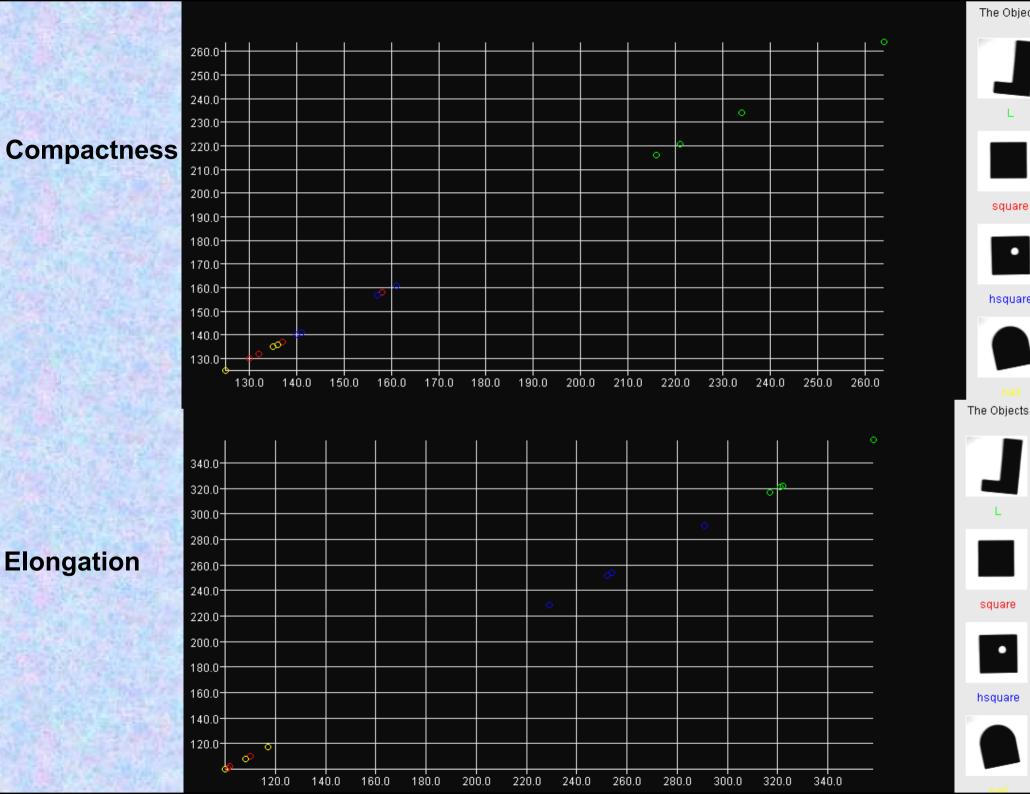
- If number of features is 2, then we get points in 2Dspace as shown in the next slide.
- We can also have an n-dimensional feature space

Decision boundary in one-dimensional case with two classes.



Decision boundary in two dimensional case with three classes





The Objects

square

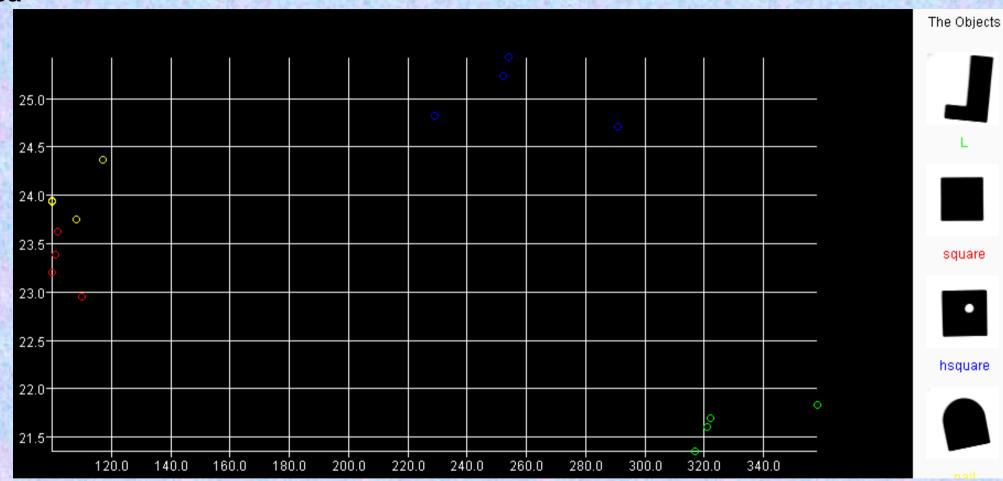
•

hsquare

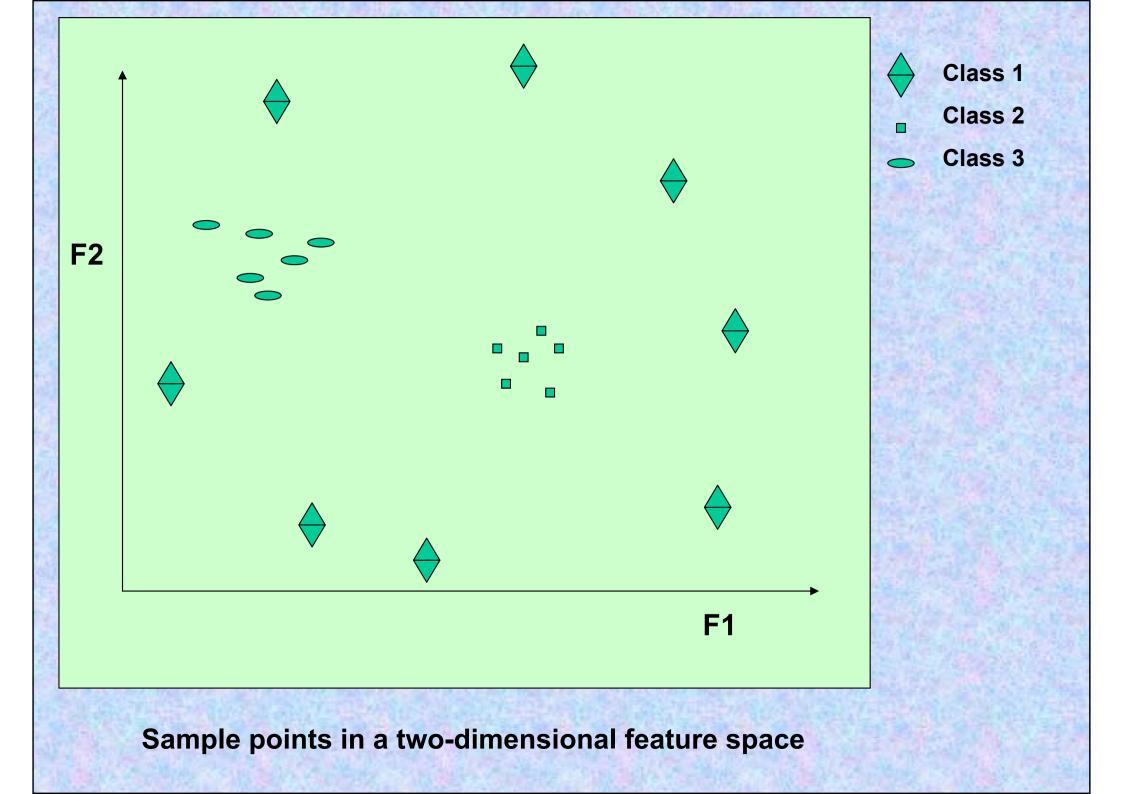
square

hsquare

Area



Elongation



Some Terminologies:

- Pattern
- Feature
- Feature vector
- Feature space
- Classification
- Decision Boundary
- Decision Region
- Discriminant function
- Hyperplanes and Hypersurfaces
- Learning
- Supervised and unsupervised
- Error
- Noise
- PDF
- Baye's Rule
- Parametric and Non-parametric approaches

Decision region and Decision Boundary

- Our goal of pattern recognition is to reach an optimal decision rule to categorize the incoming data into their respective categories
- The decision boundary separates points belonging to one class from points of other
- The decision boundary partitions the feature space into decision regions.
- The nature of the decision boundary is decided by the discriminant function which is used for decision. It is a function of the feature vector.

Hyper planes and Hyper surfaces

- For two category case, a positive value of discriminant function decides class 1 and a negative value decides the other.
- If the number of dimensions is three. Then the decision boundary will be a plane or a 3-D surface. The decision regions become semi-infinite volumes
- If the number of dimensions increases to more than three, then the decision boundary becomes a hyper-plane or a hyper-surface. The decision regions become semi-infinite hyperspaces.

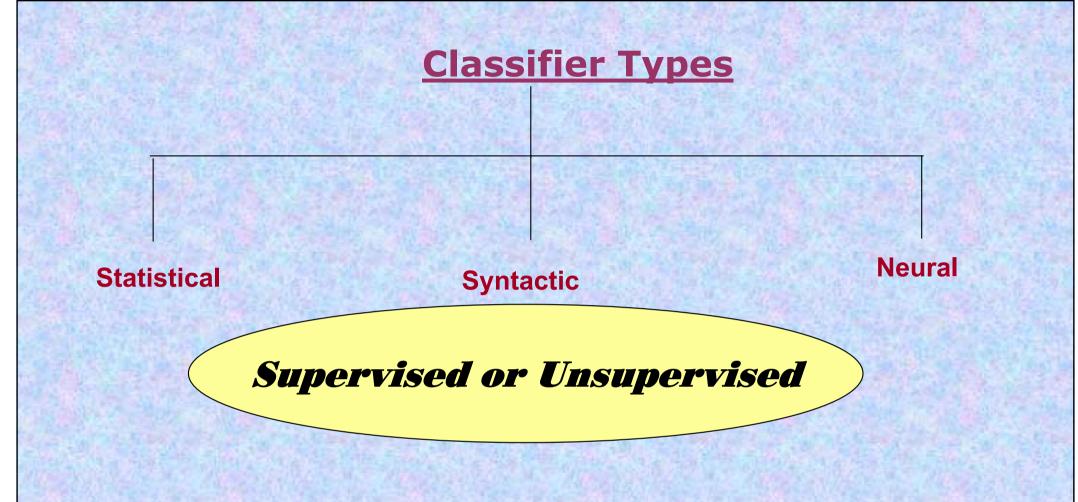
Learning

- The classifier to be designed is built using input samples which is a mixture of all the classes.
- The classifier learns how to discriminate between samples of different classes.
- If the Learning is offline i.e. Supervised method then, the classifier is first given a set of training samples and the optimal decision boundary found, and then the classification is done.
- If the learning is online then there is no teacher and no training samples (Unsupervised). The input samples are the test samples itself. The classifier learns and classifies at the same time.

Error

• The accuracy of classification depends on two things

- The optimality of decision rule used: The central task is to find an optimal decision rules which can generalize to unseen samples as well as categorize the training samples as correctly as possible. This decision theory leads to a <u>minimum error-rate classification</u>.
- The accuracy in measurements of feature vectors: This inaccuracy is because of presence of noise. Hence our classifier should deal with noisy and missing features too.



Categories of Statistical Classifiers:

- Linear
- Quadratic
- Piecewise
- Non-parametric

Parametric Decision making (Statistical) - Supervised

Goal of most classification procedures is to estimate the probabilities that a pattern to be classified belongs to various possible classes, based on the values of some feature or set of features.

In most cases, we decide which is the most likely class. We need a mathematical decision making algorithm, to obtain classification.

Bayesian decision making or Bayes Theorem

This method refers to choosing the most likely class, given the value of the feature/s. Bayes theorem calculates the probability of class membership.

Define:

- $P(w_i)$ Prior Prob. for class w_i ; P(X) Prob. for feature vector X
- P(w_i |X) Measured-conditioned or posteriori probability

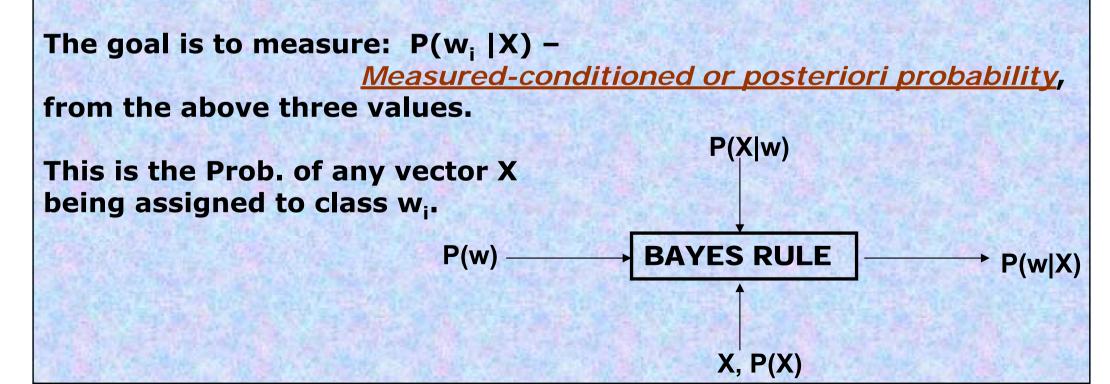
P(X | w_i) - Prob. Of feature vector X in class w_i

Bayes Theorem: $P(w_i | \vec{X}) = \frac{P(X | w_i) P(w_i)}{P(\vec{X})}$

P(X) is the probability distribution for feature X in the entire population. Also called <u>unconditional density function</u>.

 $P(w_i)$ is the <u>prior probability</u> that a random sample is a member of the class C_i .

 $P(X | w_i)$ is the <u>class conditional probability</u> of obtaining feature value X given that the sample is from class w_i . It is equal to the number of times (occurrences) of X, if it belongs to class w_i .



Take an example:

Two class problem: Cold (C) and not-cold (C'). Feature is fever (f).

Prior probability of a person having a cold, P(C) = 0.01.

Prob. of having a fever, given that a person has a cold is, P(f|C) = 0.4. Overall prob. of fever P(f) = 0.02.

Then using Bayes Th., the Prob. that a person has a cold, given that she (or he) has a fever is:

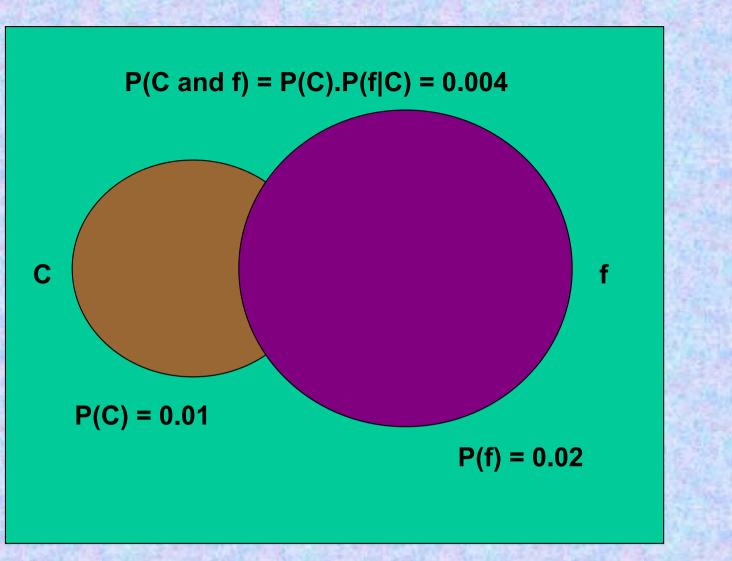
 $P(C \mid f) = \frac{P(f \mid C)P(C)}{P(f)} = \frac{0.4*0.01}{0.02} = 0.2$ Not convinced that it works? let us take an example with values to verify:

Total Population = 1000. Thus, people having cold = 10. People having both fever and cold = 4. Thus, people having only cold = 10 - 4 = 6. People having fever (with and without cold) = 0.02 * 1000 = 20. People having fever without cold = 20 - 4 = 16 (*may use this later*).

So, probability (percentage) of people having cold along with fever, out of all those having fever, is: 4/20 = 0.2 (20%).

IT WORKS, GREAT

A Venn diagram, illustrating the two class, one feature problem.



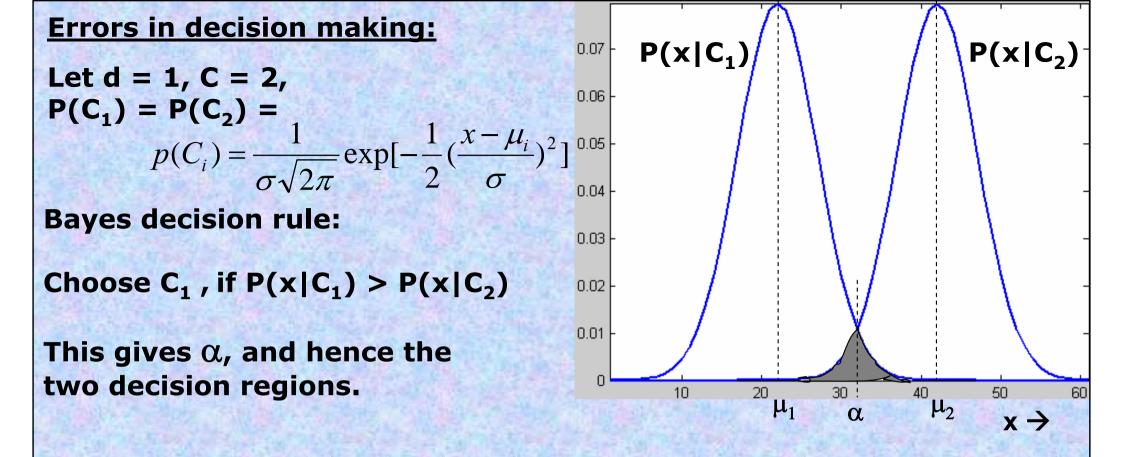
Probability of a joint event - a sample comes from class C and has the feature value X:

P(C and X) = P(C).P(X|C) = P(X).P(C|X)= 0.01*0.4 = 0.02*0.2

Also verify, for a K class problem: $P(X) = P(w_1)P(X|w_1) + P(w_2)P(X|w_2) + \dots + P(w_k)P(X|w_k)$ Thus: $P(w_i \mid \vec{X}) = \frac{P(X \mid w_i) P(w_i)}{P(w_1) P(X \mid w_1) + P(w_2) P(X \mid w_2) + \dots + P(w_k) P(X \mid w_k)}$ With our last example: P(f) = P(C)P(f|C) + P(C')P(f|C')= 0.01 * 0.4 + 0.99 * 0.01616 = 0.02

Decision or Classification algorithm according to Baye's Theorem:

Choose
$$\begin{cases} w_1; \text{ if } p(X|w_1)p(w_1) > p(X|w_2)p(w_2) \\ w_2; \text{ if } p(X|w_2)p(w_2) > p(X|w_1)p(w_1) \end{cases}$$



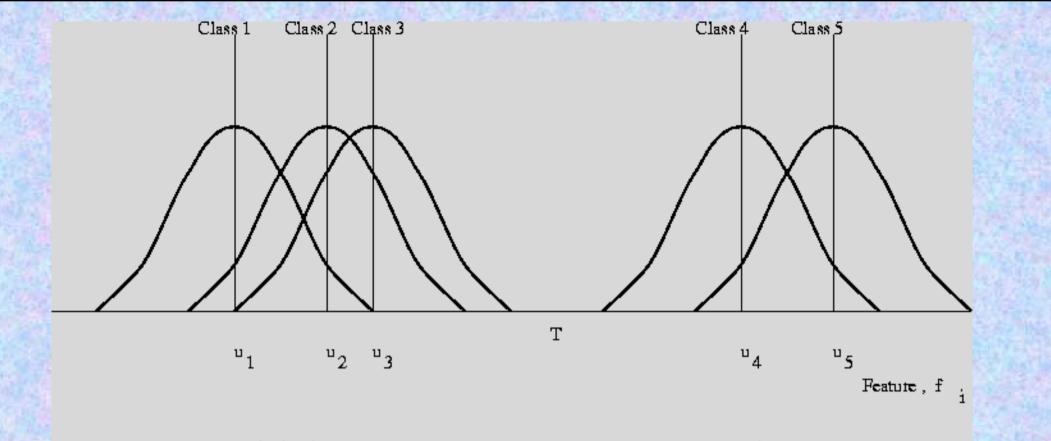
Classification error (the shaded region):

 $-\infty$

 $P(E) = P(Chosen C_1, when x belongs to C_2) + P(Chosen C_2, when x belongs to C_1)$

=
$$P(C_2) \int P(\gamma \mid C_2) d\gamma + P(C_1) \int P(\gamma \mid C_1) d\gamma$$

α



Normal distributions of feature measurement for a 5-class problem, equal variance.

A minimum distance classifier

Rule: Assign X to R_i , where X is closest to μ_i .

K-means Clustering

- Given a fixed number of k clusters, assign observations to those clusters so that the means across clusters for all variables are as different from each other as possible.
- Input
 - Number of Clusters, k
 - Collection of n, d dimensional vectors x_j, j=1, 2, ..., n
- Goal: find the k mean vectors $\mu_1, \mu_2, ..., \mu_k$
- Output
 - k x n binary membership matrix U where

$$\boldsymbol{u}_{ij} = \begin{cases} 1 & \text{if } \boldsymbol{x}_i \in \boldsymbol{G}_i \\ 0 & \text{else} \end{cases}$$

& G_j, j=1, 2, ..., k represent the k clusters

If n is the number of known patterns and c the desired number of clusters, the k-means algorithm is:

Begin

initialize n, c, $\mu_1, \mu_2, ..., \mu_c$ (randomly selected)

do 1.classify n samples according to nearest μ_i 2.recompute μ_i until no change in μ_i

return μ_1 , μ_2 , ..., μ_c

End

Classification Stage

• The samples have to be assigned to clusters in order to minimize the cost function which is:

$$J = \sum_{i=1}^{c} J_{i} = \sum_{i=1}^{c} \left| \sum_{k, x_{k} \in G_{i}} \|x_{k} - \mu_{i}\|^{2} \right|$$

- This is the Euclidian Distance of the samples from its cluster center; for all clusters this sum should be minimum
- The classification of a point x_k is done by:

$$\boldsymbol{u}_{i} = \begin{cases} 1 & \text{if } \|\boldsymbol{x}_{k} - \boldsymbol{\mu}_{i}\|^{2} \ge \|\boldsymbol{x}_{k} - \boldsymbol{\mu}_{j}\|^{2} , \forall k \neq i \\ 0 & \text{otherwise} \end{cases}$$

Re-computing the Means

The means are recomputed according to:

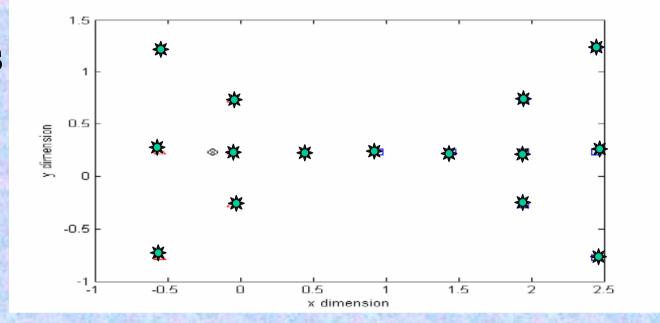
$$\mu_i = \frac{1}{|G_i|} \left(\sum_{k, x_k \in G_i} x_k \right)$$

Disadvantages

- What happens when there is overlap between classes... that is a point is equally close to two cluster centers..... Algorithm will not terminate
- The Terminating condition is modified to "Change in cost function (computed at the end of the Classification) is below some threshold rather than 0".

An Example

- The no of clusters is two in this case.
- But still there is some overlap



	Membership Matrix U														
Point s(k)	1	2	3	4	5	6	7	8	9	10	11	12	13	14	15
u _{1k}	1	1	1	1	1	1	1	0	0	0	0	0	0	0	0
u _{2k}	0	0	0	0	0	0	0	1	1	1	1	1	1	1	1

Some necessary elements of

Probability theory and Statistics

Normal Density: $p(x) = \frac{1}{\sigma\sqrt{2\pi}} \exp\left[-\frac{1}{2}\left(\frac{x-\mu}{\sigma}\right)^2\right]$

Bivariate Normal Density:

$$p(x,y) = \frac{e^{-\frac{1}{2(1-\rho_{xy}^2)}\left[(\frac{x-\mu_x}{\sigma_x})^2 - \frac{2\rho_{xy}(x-\mu_x)(y-\mu_y)}{\sigma_x\sigma_y} + (\frac{y-\mu_y}{\sigma_y})^2\right]}}{2\pi\sigma_x\sigma_y\sqrt{(1-\rho_{xy}^2)}}$$

 μ - Mean; σ - S.D.; ρ_{xy} - Correlation Coefficient

Visualize ρ as equivalent to the orientation of the 2-D Gabor filter.

For x as a discrete random variable, the expected value of x:

$$E(x) = \sum_{i=1}^{n} x_i P(x_i) = \mu_x$$

E(x) is also called the first moment of the distribution. The kth moment is defined as: $E(x^k) = \sum_{i=1}^{n} x_i^k P(x_i)$

 $P(x_i)$ is the probability of $x = x_i$.

Second, third,... moments of the distribution p(x) are the expected values of: x^2 , x^3 ,...

The kth central moment is defined as:

$$E[(x - \mu_x)^k] = \sum_{i=1}^{m} (x - \mu_x)^k P(x_i)$$

Thus, the second central moment (also called Variance) of a random variable x is defined as:

$$\sigma_x^2 = E[\{x - E(x)\}^2] = E[(x - \mu_x)^2]$$

S.D. of x is σ_x .

$$\sigma_x^2 = E[\{x - E(x)\}^2] = E[(x - \mu_x)^2]$$
$$= E(x^2) - 2\mu_x^2 + \mu_x^2 = E(x^2) - \mu_x^2$$

Thus

$$E(x^2) = \sigma^2 + \mu^2$$

If z is a new variable: z = ax + by; Then E(z) = E(ax + by)=aE(x) + bE(y).

Covariance of x and y, is defined as: $\sigma_{xy} = E[(x - \mu_x)(y - \mu_y)]$

Covariance indicates how much x and y vary together. The value depends on how much each variable tends to deviate from its mean, and also depends on the degree of association between x and y.

Correlation between x and y:
$$\rho_{xy} = \frac{\sigma_{xy}}{\sigma_x \sigma_y} = E[(\frac{x - \mu_x}{\sigma_x})(\frac{y - \mu_y}{\sigma_y})]$$

Property of correlation coefficient: $-1 \le \rho_{xy} \le 1$

For Z:

$$E[(z - \mu_z)^2] = a^2 \sigma_x^2 + 2ab \sigma_{xy} + b^2 \sigma_y^2$$

If $\sigma_{xy} = 0$, $\sigma_z^2 = a^2 \sigma_x^2 + b^2 \sigma_y^2$

1.00.8 0.40.0-0.8-1.0-0.41.01.01.00.0 -1.0-1.0-1.00.00.00.00.00.00.00.0 $\sigma_{xy} =$ $E[(\frac{x-\mu_x}{x})(\frac{y-\mu_y}{x})]$ ρ_{xy} $\sigma_{x} = \frac{\sigma_{y} E(XY) - E(X)E(Y)}{\sqrt{E(X^{2}) - E^{2}(X)} \sqrt{E(Y^{2}) - E^{2}(Y)}}$ $\sigma_r \sigma_1$

The <u>correlation coefficient</u> can also be viewed as the cosine of the angle between the two vectors (R^{D}) of samples drawn from the two random variables.

This method only works with centered data, i.e., data which have been shifted by the sample mean so as to have an average of zero.

<u>Multi-variate Case</u>: $X = [\mathbf{x}_1 \ \mathbf{x}_2 \ \dots \ \mathbf{x}_d]^T$

Mean vector:

 μ_1

 μ_2

 μ_d

Covariance matrix (symmetric):

d-dimensional normal density is:

$$p(X) = \frac{1}{\sqrt{\det(\Sigma)(2\pi)^d}} \exp\left[-\frac{(X-\mu)^T \Sigma^{-1}(X-\mu)}{2}\right]$$
$$= \frac{1}{\sqrt{\det(\Sigma)(2\pi)^d}} \exp\left[-\frac{1}{2} \sum_{ij} (x_i - \mu_i) s_{ij} (x_j - \mu_j)\right]$$

$$p(X) = \frac{1}{\sqrt{\det(\Sigma)(2\pi)^{d}}} \exp\left[-\frac{(X-\mu)^{T}\Sigma^{-1}(X-\mu)}{2}\right]$$

$$= \frac{1}{\sqrt{\det(\Sigma)(2\pi)^{d}}} \exp\left[-\frac{1}{2}\sum_{ij}(x_{i}-\mu_{i})s_{ij}(x_{j}-\mu_{j})\right]$$
where s_{ij} is the ijth component of Σ^{-1} (the inverse of covariance matrix Σ).
Special case, d = 2; where X = (x y)^T; Then:
and

$$\Sigma = \begin{pmatrix} \sigma_{x}^{2} & \sigma_{xy} \\ \sigma_{xy} & \sigma_{y}^{2} \end{pmatrix} = \begin{pmatrix} \sigma_{x}^{2} & \rho_{xy}\sigma_{x}\sigma_{y} \\ \rho_{xy}\sigma_{x}\sigma_{y} & \sigma_{y}^{2} \end{pmatrix}$$
Can you now obtain this,
as given earlier:

$$p(x, y) = \frac{e^{-\frac{1}{2(1-\rho_{xy}^{2})}\left[(\frac{x-\mu_{x}}{\sigma_{x}})^{2} - \frac{2\rho_{xy}(x-\mu_{x})(y-\mu_{y})}{\sigma_{x}\sigma_{y}} + (\frac{y-\mu_{y}}{\sigma_{y}})^{2}}\right]}{2\pi\sigma_{x}\sigma_{y}\sqrt{(1-\rho_{xy}^{2})}}$$

Sample mean is defined as:

$$\tilde{x} = \sum_{i=1}^{n} x_i P(x_i) = \frac{1}{n} \sum_{i=1}^{n} x_i$$

where, $P(x_i) = 1/n$.

Sample Variance is:
$$\sigma_x^2 = \frac{1}{n} \sum_{i=1}^n (x_i - x)^2$$

Higher order moments may also be computed:

$$E(x_i - x)^3; E(x_i - x)^2$$

Covariance of a bivariate distribution:

$$\sigma_{xy} = E[(x - \mu_x)(y - \mu_y)] = \frac{1}{n} \sum_{i=1}^n (x - x)(y - y)$$

MAXIMUM LIKELIHOOD ESTIMATE

The ML estimate (MLE) of a parameter is that value which, when substituted into the probability distribution (or density), produces that distribution for which the probability of obtaining the entire observed set of samples is maximized.

Problem: Find the maximum likelihood estimate for μ in a normal distribution.

Normal Density:

$$p(x) = \frac{1}{\sigma\sqrt{2\pi}} \exp\left[-\frac{1}{2}\left(\frac{x-\mu}{\sigma}\right)^2\right]$$

Assuming all random samples to be independent:

$$p(x_1, ..., x_n) = p(x_1)....p(x_n) = \prod_{i=1}^n p(x_i)$$
$$= \frac{1}{\sigma^n (2\pi)^{n/2}} \exp[-\frac{1}{2\sigma^2} \sum_{i=1}^n (\frac{x-\mu}{\sigma})^2]$$

Taking derivative (w.r.t. μ) of the LOG of the above:

$$\frac{1}{2\sigma^2} \sum_{i=1}^n (x_i - \mu) \cdot 2 = \frac{1}{\sigma^2} \left[\sum_{i=1}^n x_i - n\mu \right]$$

Setting this term = 0, we get:

$$\mu = \frac{1}{n} \sum_{i=1}^{n} x_i = \tilde{x}$$

Decision Regions and Boundaries

A classifier partitions a feature space into class-labeled <u>decision regions (DRs)</u>.

If decision regions are used for a possible and unique class assignment, the regions must cover R^d and be disjoint (nonoverlapping. In Fuzzy theory, decision regions may be overlapping.

The border of each decision region is a **Decision Boundary (DBs)**.

Typical classification approach is as follows:

Determine the decision region (in R^d) into which X falls, and assign X to this class.

This strategy is simple. But determining the DRs is a challenge.

It may not be possible to visualize, DRs and DBs, in a general classification task with a large number of classes and higher feature space (dimension).

Classifiers are based on *Discriminant functions*.

In a C-class case, Discriminant functions are denoted by: g_i(X), i = 1,2,...,C.

This partitions the R^d into C distinct (disjoint) regions, and the process of classification is implemented using the <u>Decision Rule</u>:

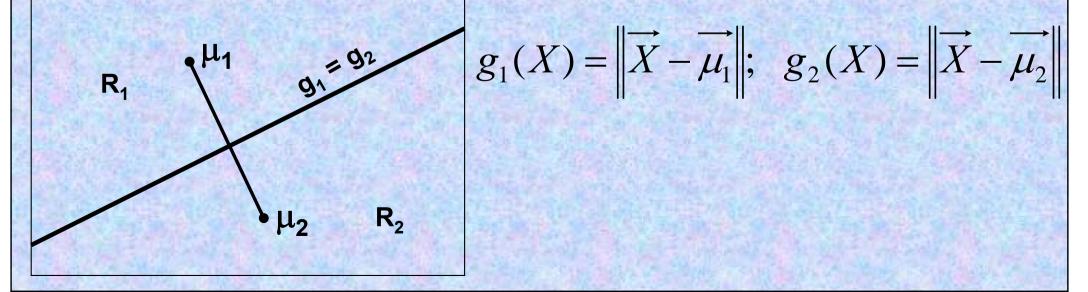
Assign X to class C_m (or region m), where: $g_m(X) > g_i(X), \forall i, i \neq m$.

Decision Boundary is defined by the locus of points, where:

$$g_k(X) = g_l(X), k \neq l$$

Minimum distance (also NN) classifier:

Discriminant function is based on the distance to the class mean:



Let the discrimination function for the ith class be: $g_i(X) = P(C_i | X)$, and assume $P(C_i) = P(C_i), \forall i, j; i \neq j$. **Remember, multivariate Gaussian density?** $g_{i}(X) = P(X | C_{i}) = \frac{1}{\sqrt{\det(\Sigma)(2\pi)^{d}}} \exp[-\frac{(X - \mu_{i})^{T} \Sigma_{i}^{-1} (X - \mu_{i})}{2}]$ Define: Define: $G_{i}(X) = \log[P(X | C_{i})] = \log[\frac{1}{\sqrt{\det(\Sigma)(2\pi)^{d}}}] - \frac{(X - \mu_{i})^{T} \Sigma_{i}^{-1} (X - \mu_{i})}{2}$ $= k \cdot \vec{d}_i^2 + q$

Thus the classification is now influenced by the square distance (hyper-dimensional) of X from μ_i , weighted by the Σ^{-1} . Let us examine: $\vec{d}_i^2 = (X - \mu_i)^T \Sigma_i^{-1} (X - \mu_i)$

This quadratic term (scalar) is known as the <u>Mahalanobis distance</u> (the distance from X to μ_i in feature space).

$$\vec{d}_{i}^{2} = (X - \mu_{i})^{T} \Sigma_{i}^{-1} (X - \mu_{i})$$

For a given X, some $G_m(X)$ is largest and also $(d_m)^2$ is the smallest, for a class i = m (assign X to class m, based on NN Rule).

Simplest case: $\Sigma = I$, the criteria becomes the Euclidean distance norm.

This is equivalent to obtaining the mean μ_m , for which X is the nearest, for all μ_i . The distance function is then:

$$\vec{d}_{i}^{2} = \|X - \mu_{i}\|^{2} = X^{T}X - 2\mu_{i}^{T}X + \mu_{i}^{T}\mu \quad \text{(all vector notations)}$$

Thus, $G_{i}(X) = \frac{d_{i}^{2}}{2} = \frac{(X^{T}X)}{2} - \frac{\mu_{i}^{T}X}{4} + \frac{(\mu_{i}^{T}\mu_{i})}{2}$

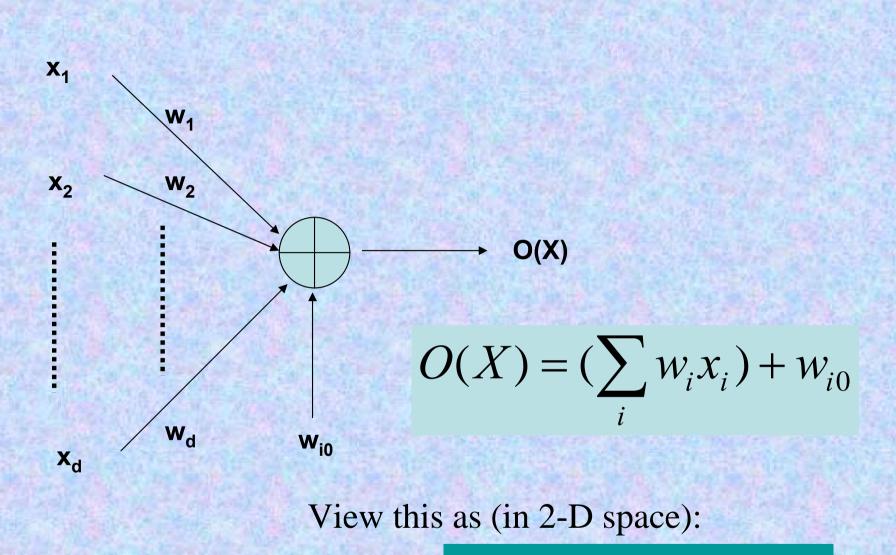
 $= \omega_i^T X + \omega_{i0}$

Neglecting the class-invariant term.

where, $\omega_i^T = \mu_i$ and $\omega_{i0} = -\frac{\mu_i^T \mu_i}{2}$

This gives the simplest linear discriminant function or correlation detector.

The perceptron (ANN) built to form the linear discriminant function



Y = MX + C

The decision region boundaries are determined by solving :

$$G_i(X) = G_j(X)$$
, which gives : $(\omega_i^T - \omega_j^T)X + (\omega_{i0} - \omega_{j0}) = 0$

This is an expression of a hyperplane separating the decision regions in R^d. The hyperplane will pass through the origin, if: $\omega_{i0} = \omega_{i0}$

Generalized results (Gaussian case) of a discriminant function:

$$G_{i}(X) = \log[P(X | C_{i})] = \log[\frac{1}{\sqrt{\det(\Sigma_{i})(2\pi)^{d}}}] - \frac{(X - \mu_{i})^{T} \Sigma_{i}^{-1} (X - \mu_{i})}{2}$$
$$= -\frac{1}{2} (X - \mu_{i})^{T} \Sigma_{i}^{-1} (X - \mu_{i}) - (\frac{d}{2}) \log(2\pi) - \frac{1}{2} \log(\Sigma_{i})$$

The mahalanobis distance (quadratic term) spawns a number of different surfaces, depending on Σ^{-1} . It is basically a vector distance using a Σ^{-1} norm. It is denoted as: $\|X - \mu_i\|_{\Sigma^{-1}}^2$

Make the case of Baye's rule more general for class assignment. Earlier we has assumed that: $g_i(X) = P(C_i | X)$, assuming $P(C_i) = P(C_i), \forall i, j; i \neq j$. Now, $G_i(X) = \log[P(C_i | X) \cdot P(X)] = \log[P(X | C_i)] + \log[P(C_i)]$ $G_{i}(X) = \log[\frac{1}{\sqrt{\det(\Sigma_{i})(2\pi)^{d}}}] - \frac{(X - \mu_{i})^{T} \Sigma_{i}^{-1} (X - \mu_{i})}{2} + \log[P(C_{i})]$ $= -\frac{1}{2}(X - \mu_i)^T \Sigma_i^{-1} (X - \mu_i) - (\frac{d}{2})\log(2\pi) - \frac{1}{2}\log(\Sigma_i) + \log[P(C_i)]$ $= -\frac{1}{2} (X - \mu_i)^T \Sigma_i^{-1} (X - \mu_i) - \frac{1}{2} \log(\Sigma_i) + \log[P(C_i)]$ Neglecting the constant term <u>Simpler case</u>: $\Sigma_i = \sigma^2 I$, and eliminating the class-independent bias, we have: $G_{i}(X) = -\frac{1}{2\sigma^{2}}(X - \mu_{i})^{T}(X - \mu_{i}) + \log[P(C_{i})]$

These are loci of constant hyper-spheres, centered at class mean.

If Σ is a diagonal matrix, with equal/unequal σ_{ii}^2 : onal matrix, with equal/unequal $\sigma_{ii}^{(1)}$ $\Sigma = \begin{bmatrix} \sigma_1^2 & 0 & \dots & 0 \\ 0 & \sigma_2^2 & \dots & 0 \\ \dots & \dots & \dots & \dots \\ 0 & 0 & \dots & \sigma_d^2 \end{bmatrix} and \Sigma^{-1} = \begin{bmatrix} 1/\sigma_1^2 & 0 & \dots & 0 \\ 0 & 1/\sigma_2^2 & \dots & 0 \\ \dots & 0 & 0 & \dots & 0 \\ \dots & \dots & \dots & \dots \\ 0 & 0 & \dots & 1/\sigma_d^2 \end{bmatrix}$

Considering the discriminant function:

$$G_i(X) = -\frac{1}{2}(X - \mu_i)^T \Sigma_i^{-1}(X - \mu_i) - \frac{1}{2}\log(\Sigma_i) + \log[P(C_i)]$$

This now will yield a weighted distance classifier. Depending on the covariance term (more spread/scatter or not), we tend to put more emphasis on some feature vector components than the other.

Check out the following:

This will give hyper-elliptical surfaces in R^d, for each class.

It is also possible to linearise it.

More general decision boundaries

Take $P(C_i) = K$ for all i, and eliminating the class independent terms yield:

$$G_{i}(X) = (X - \mu_{i})^{T} \Sigma_{i}^{-1} (X - \mu_{$$

 $\vec{d}_{i}^{2} = (X - \mu_{i})^{T} \Sigma_{i}^{-1} (X - \mu_{i}) = -X^{T} \Sigma_{i}^{-1} X + 2\mu_{i}^{T} \Sigma_{i}^{-1} X - \mu_{i}^{T} \Sigma_{i}^{-1} \mu_{i}$

 $G_i(X) = (\Sigma_i^{-1}\mu_i)^T X - \frac{1}{2}\mu_i^T \Sigma_i^{-1}\mu_i \quad \text{as } \Sigma_i = \Sigma, \text{ and are symmetric.}$

Thus, $G_i(X) = \omega_i^T X + \omega_{i0}$

where
$$\omega_i = \Sigma_i^{-1} \mu_i$$
 and $\omega_{i0} = -\frac{1}{2} \mu_i^T \Sigma_i^{-1} \mu_i$

L

Thus the decision surfaces are hyperplanes and decision boundaries will also be linear (use $G_i(X) = G_i(X)$, as done earlier)

Beyond this, if a diagonal Σ is class-dependent or off-diagonal terms are non-zero, we get non-linear DFs, DRs or DBs.

The discriminant function (DF) for linearly separable classes is: $g_i(X) = \omega_i^T X + \omega_{i0}$

where, ω_i is a dx1 vector of weights used for class i.

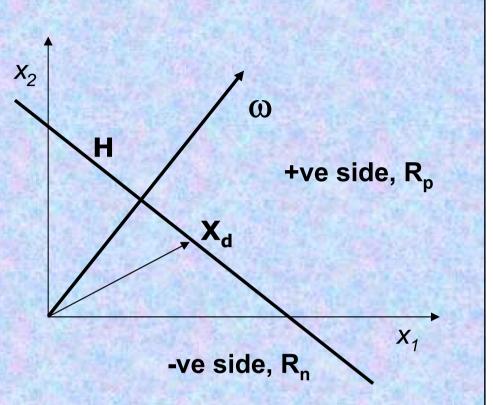
This function leads to DBs that are hyperplanes. It's a point in 1D, line in 2-D, planar surfaces in 3-D, and

3-D case: $(\omega_1 \omega_2 \omega_3) \begin{pmatrix} x_1 \\ x_2 \\ x_3 \end{pmatrix} = 0$ is a plane passing through the origin. In general, the equation: $\omega^T (\vec{X} - \vec{X}_d) = 0; => \omega^T \vec{X} - d = 0$ represents a plane H passing through any point (position vector) X_{d} .

This plane partitions the space into two mutually exclusive regions, $\omega^{T} \overrightarrow{X} - d \begin{cases} > 0 & \text{if } X \in R_{p} \\ = 0 & \text{if } X \in H \\ < 0 & \text{if } X \in R_{n} \end{cases}$ say R_p and R_n. The assignment of the vector X to either the +ve side, or -ve side or along H, can be implemented by:

Linear Discriminant Function g(X):

$$g(X) = \omega^T \overline{X} - d$$



Orientation of H is determined by ω .

Location of H is determined by d.

H is a hyperplane for d > 3. The figure shows a 2D representation.

Quadratic Decision Boundaries

In \mathbb{R}^d with $X = (x_1, x_2, ..., x_d)^T$, consider the equation:

$$\sum_{i=1}^{d} w_{ii} x_i^2 + \sum_{i=1}^{d-1} \sum_{j=i+1}^{d} w_{ij} x_i x_j + \sum_{i=1}^{d} w_i x_i + w_o = 0 \quad ..1$$

The above equation is defined by a quadric discriminant function, which yields a quadric surface.

If d=2, $X = (x_1, x_2)^T$ equation (1) becomes:

 $w_{11}x_1^2 + w_{22}x_2^2 + w_{12}x_1x_2 + w_1x_1 + w_2x_2 + w_0 = 0 \quad ..2$

Special cases of equation:

 $w_{11}x_1^2 + w_{22}x_2^2 + w_{12}x_1x_2 + w_1x_1 + w_2x_2 + w_0 = 0 \quad ..2$

Case 1:

 $w_{11} = w_{22} = w_{12} = 0$; Eqn. (2) defines a line.

Case 2:

 $w_{11} = w_{22} = K; w_{12} = 0;$ defines a circle.

Case 3:

 $w_{11} = w_{22} = 1$; $w_{12} = w_1 = w_2 = 0$; defines a circle whose center is at the origin.

Case 4:

 $w_{11} = w_{22} = 0$; defines a bilinear constraint.

Case 5:

 $w_{11} = w_{12} = w_2 = 0$; defines a parabola with a specific orientation.

Case 6:

: $w_{11} \neq 0, w_{22} \neq 0, w_{11} \neq w_{22}; w_{12} = w_1 = w_2 = 0$ defines a simple ellipse.

Selecting suitable values of w_i 's, gives other conic sections. For d \geq 3, we define a family of hyper-surfaces in R^d.

$$\sum_{i=1}^{d} w_{ii} x_i^2 + \sum_{i=1}^{d-1} \sum_{j=i+1}^{d} w_{ij} x_i x_j + \sum_{i=1}^{d} w_i x_i + \omega_o = 0 \quad ..1$$

In the above equation, the number of parameters:

2d + 1 + d(d-1)/2 = (d+1)(d+2)/2.

Organize these parameters, and manipulate the equation to obtain:

$$X^{T}WX + w^{T}X + \omega_{o} = 0 \qquad ..3$$

w has d terms, ω_{o} has one term, and W (ω_{ij}) is a dxd matrix as:

(d²-d) non-diagonal terms of the matrix W, is obtained by duplicating (split into two parts): $\omega_{ij} = \begin{cases} w_{ii} & \text{if } i = j \\ \frac{1}{2}w_{ij} & \text{if } i \neq j \\ \frac{1}{2}w_{ij} & \text{if } i \neq j \end{cases}$

In equation 3, the symmetric part of matrix W, contributes to the Quadratic terms. Equation 3 generally defines a hyperhyperboloidal surface. If W = I, we get a hyperplane.

$$\overline{X}^T W \overline{X} + w^T \overline{X} + \omega_o = 0$$

$$\vec{d}_i^2 = (X - \mu_i)^T \Sigma^{-1} (X - \mu_i) = -X^T \Sigma^{-1} X + 2\mu_i^T \Sigma^{-1} X - \mu_i^T \Sigma^{-1} \mu_i$$

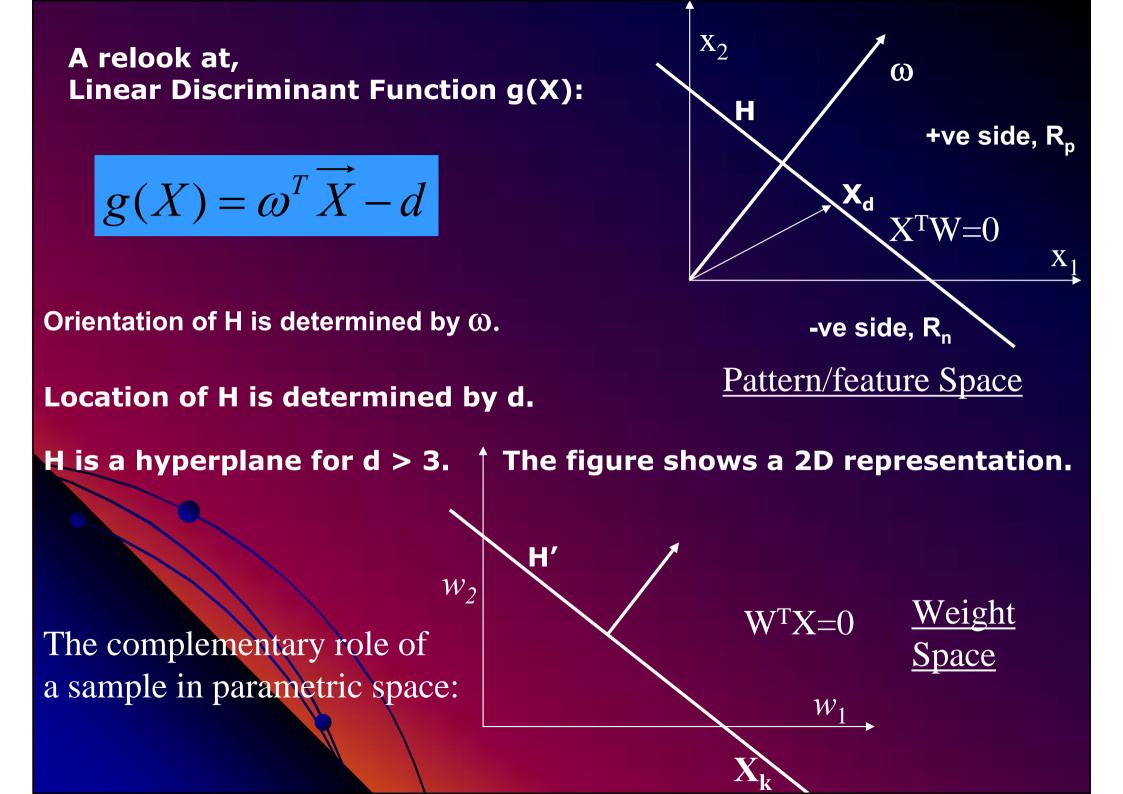
Example of linearization:

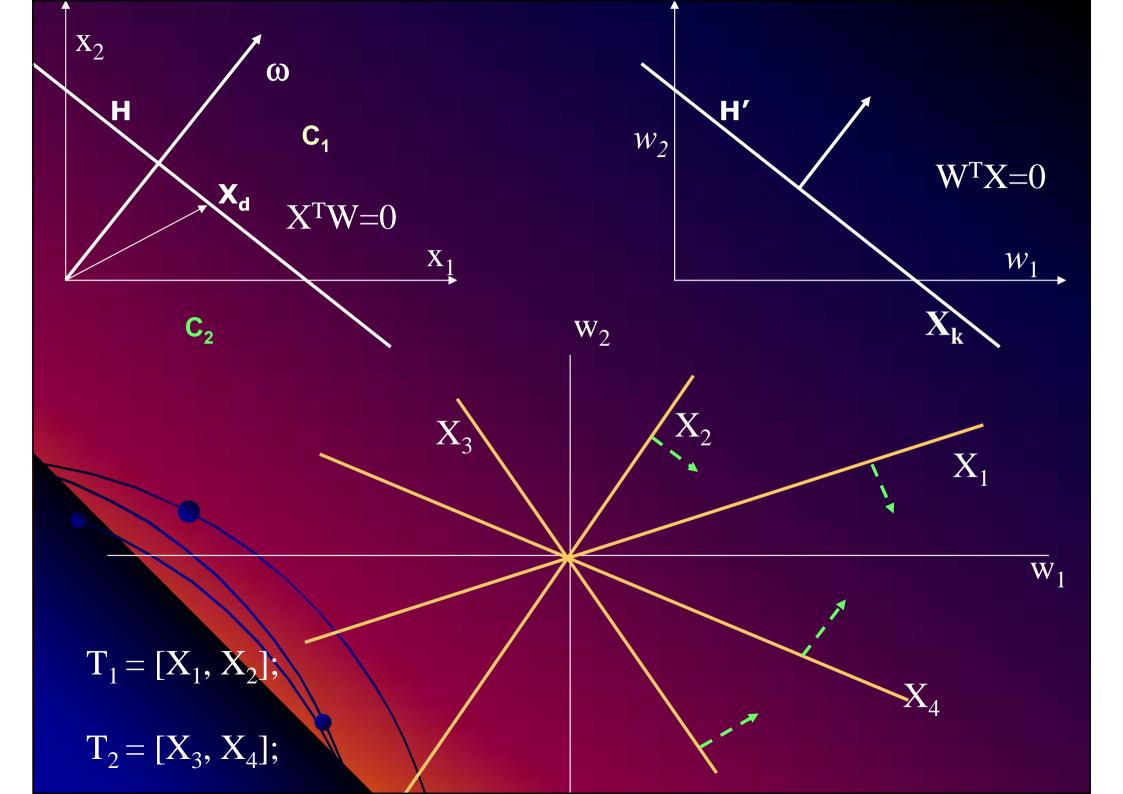
$$g(X) = x_2 - x_1^2 - 3x_1 + 6 = 0$$

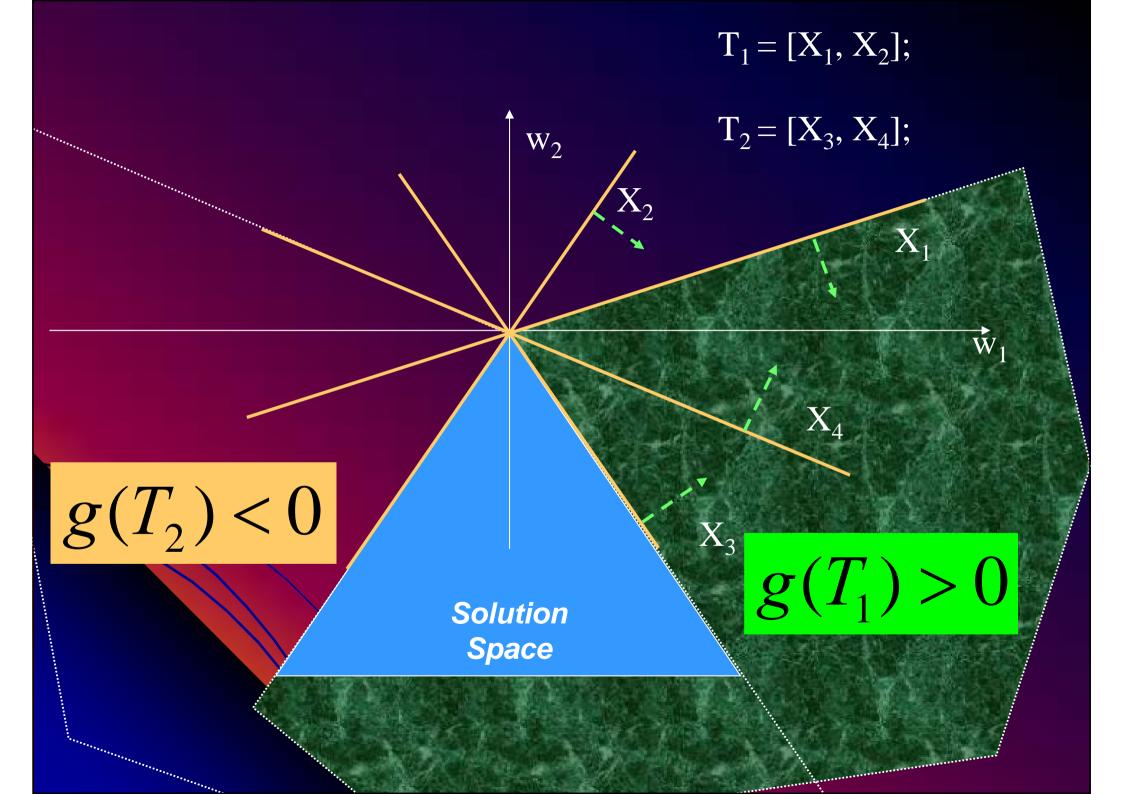
To **Linearize**, let $x_3 = x_1^2$. Then:

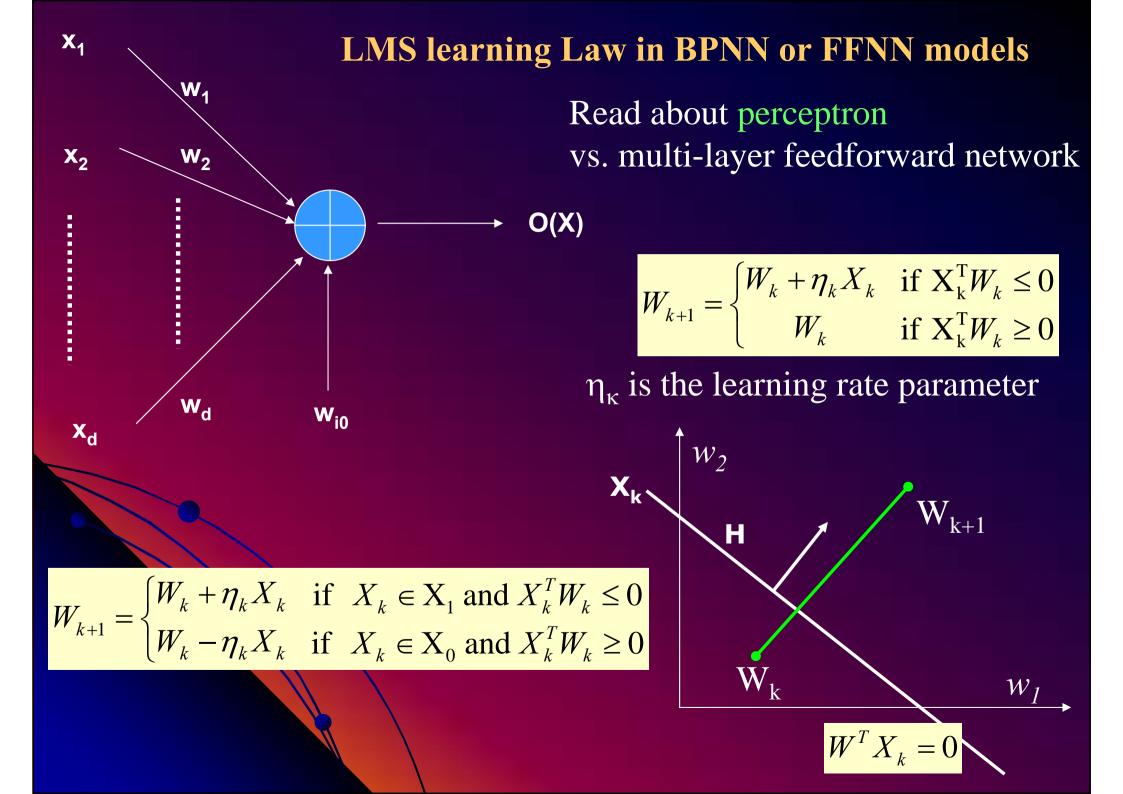
$$g(X) = x_2 - x_3 - 3x_1 + 6 = W^T X + w_o$$

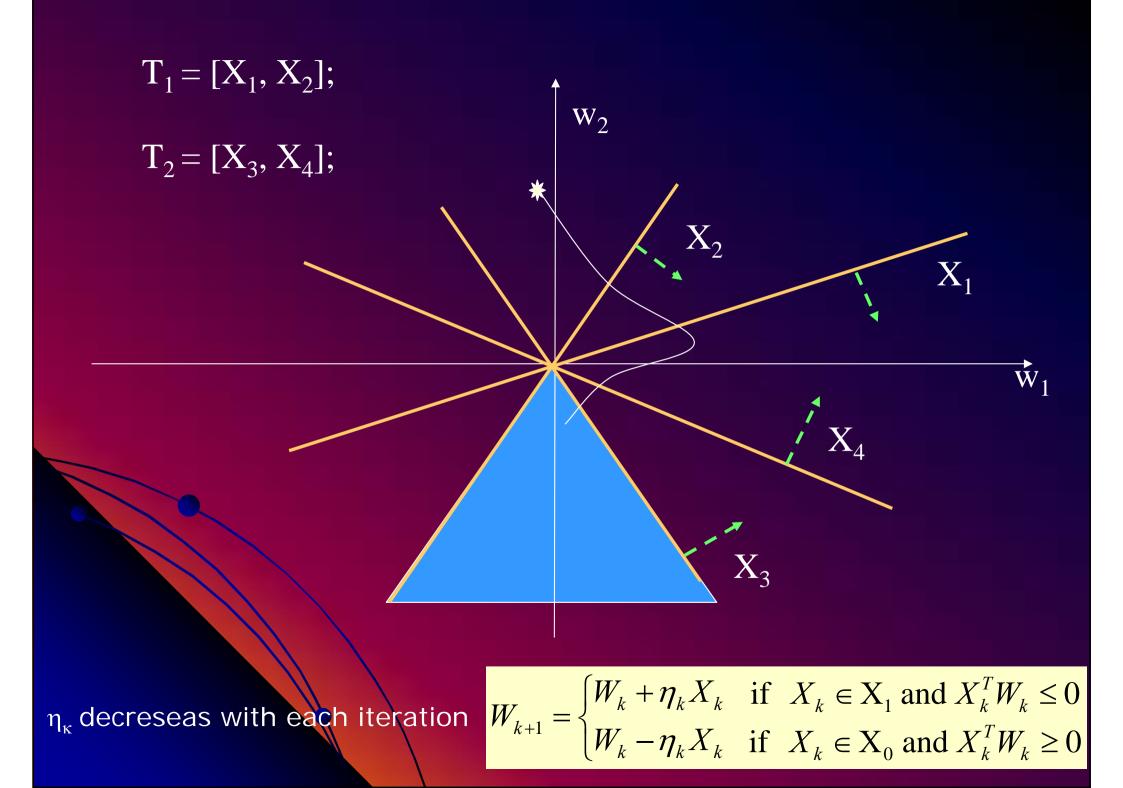
where, $X = [x_1, x_2, x_3]^T$
and $W^T = [-3, 1, -1]$



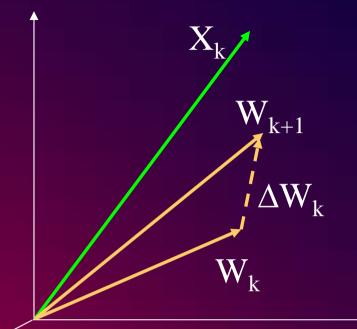








In case of FFNN, the objective is to minimize the error term:



$$e_k = d_k - s_k = d_k - X_k^T W_k$$

$$\alpha - LMS$$
 Learning Algorithm :
 $\Delta W_k = \eta e_k \stackrel{\uparrow}{X_k}$

Principal Component Analysis

Eigen analysis, Karhunen-Loeve transform

Eigenvectors: derived from Eigen decomposition of the scatter matrix

A projection set that best explains the distribution of the representative features of an object of interest.

PCA techniques choose a dimensionality-reducing linear projection that maximizes the scatter of all projected samples.

Principal Component Analysis Contd.

• Let us consider a set of N sample images $\{x_1, x_2, \dots, x_N\}$ taking values in *n*-dimensional image space.

• Each image belongs to one of c classes $\{X_1, X_2, \dots, X_c\}$.

• Let us consider a linear transformation, mapping the original *n*-dimensional *image space* to *m*-dimensional *feature space*, where m < n.

• The new feature vectors $y_k \in \mathbb{R}^m$ are defined by the linear transformation –

$$y_k = W^T x_k \qquad k = 1, 2, \dots, N$$

where $W \in \mathbb{R}^{n \times m}$ is a matrix with orthogonal columns representing the basis in feature space.

Principal Component Analysis Contd..

• Total scatter matrix S_T is defined as

$$S_T = \sum_{k=1}^{N} (x_k - \mu) (x_k - \mu)^T$$

where, *N* is the number of samples , and $\mu \in \mathbb{R}^n$ is the mean image of all samples .

• The scatter of transformed feature vectors $\{y_1, y_2, \dots, y_N\}$ is $W^T S_T W$.

• In PCA, W_{opt} is chosen to maximize the determinant of the total scatter matrix of projected samples, *i.e.*,

$$W_{opt} = \arg\max_{W} |W^T S_T W|$$

where $\{w_i \mid i = 1, 2, ..., m\}$ is the set of *n* dimensional eigenvectors of S_T corresponding to *m* largest eigenvalues.

Principal Component Analysis Contd.

• Eigenvectors are called eigen images/pictures and also basis images/facial basis for faces.

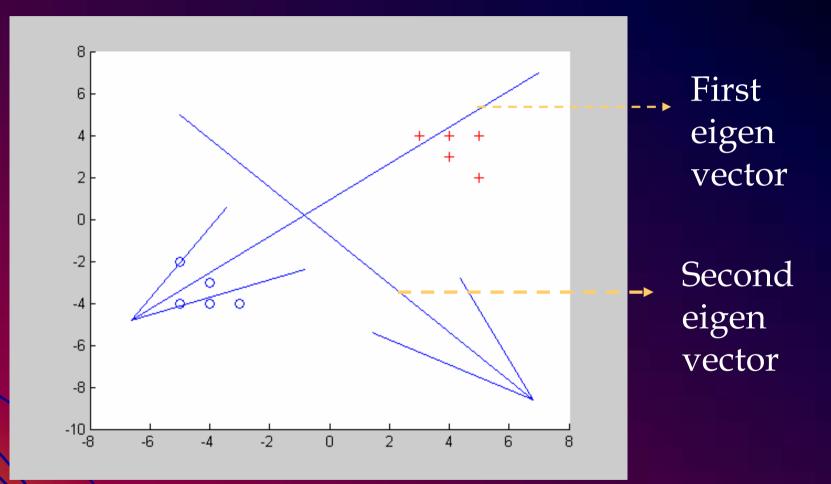
• Any face can be reconstructed approximately as a weighted sum of a small collection of images that define a facial basis (eigen images) and a mean image of the face.

 Data form a scatter in the feature space through projection set (eigen vector set)

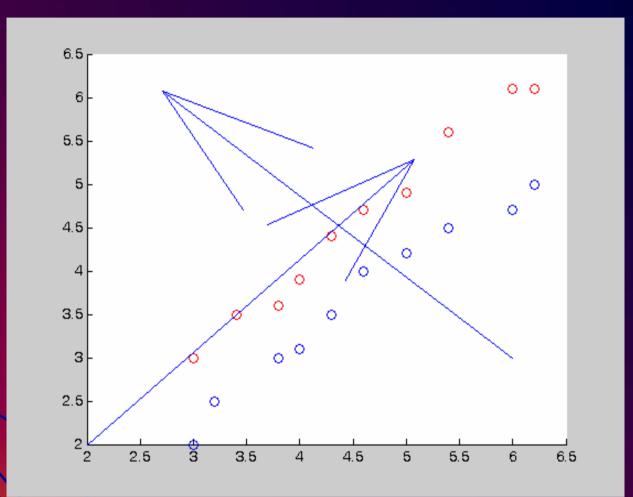
 Features (eigenvectors) are extracted from the training set without prior class information

→Unsupervised learning

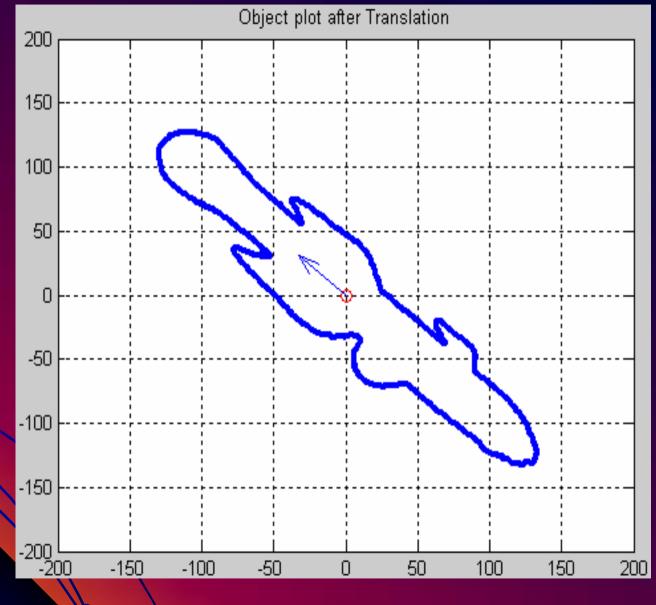
Demonstration of KL Transform



Another One



Another Example



Source: SQUID Homepage

Principal components analysis (PCA) is a technique used to reduce multi-dimensional <u>data sets</u> to lower dimensions for analysis.

The applications include <u>exploratory data analysis</u> and generating predictive models. PCA involves the computation of the <u>eigenvalue decomposition</u> or <u>Singular value decomposition</u> of a data set, usually after mean centering the data for each attribute.

PCA is mathematically defined as an <u>orthogonal linear</u> <u>transformation</u>, that transforms the data to a new <u>coordinate</u> <u>system</u> such that the <u>greatest variance</u> by any projection of the data comes to lie on the <u>first coordinate</u> (called the first principal component), the second greatest variance on the second coordinate, and so on.

PCA can be used for <u>dimensionality reduction</u> in a data set by retaining those characteristics of the data set that contribute most to its <u>variance</u>, by keeping lower-order principal components and ignoring higher-order ones. Such low-order components often contain the "most important" aspects of the data. But this is not necessarily the case, depending on the application. For a data <u>matrix</u>, X^T, with zero <u>empirical mean</u> (the empirical mean of the distribution has been subtracted from the data set), where each *column* is made up of results for a different subject, and each *row* the results from a different probe. This will mean that the PCA for our data matrix X will be given by:

 $Y = W^T X = \Sigma V,$

where $W\Sigma V^T$ is the singular value decomposition (SVD) of X.

Goal of PCA:

Find some orthonormal matrix W^T, where Y = W^TX; such that

 $COV(Y) \equiv (1/(n-1))YY^T$ is diagonalized.

The rows of W are the principal components of X, which are also the eigenvectors of COV(X).

Unlike other linear transforms (DCT, DFT, DWT etc.), PCA does not have a fixed set of <u>basis vectors</u>. Its basis vectors depend on the data set. The Karhunen-Loève transform is therefore equivalent to finding the <u>singular value decomposition</u> of the data matrix X, and then obtaining the reduced-space data matrix Y by projecting X down into the reduced space defined by only the first L singular vectors, W_L :

$$X = W \Sigma V^T; \quad Y = W_L^T X = \Sigma_L V_L^T$$

The matrix W of singular vectors of X is equivalently the matrix W of eigenvectors of the matrix of observed covariances $C = X X^{T}$ (find out?) =:

$$COV(X) = XX^T = W\Sigma\Sigma^T W^T = WDW^T$$

The <u>eigenvectors</u> with the largest <u>eigenvalues</u> correspond to the dimensions that have the strongest <u>correlation</u> in the data set. PCA is equivalent to <u>empirical</u> <u>orthogonal functions</u> (EOF).

PCA is a popular technique in <u>pattern recognition</u>. But it is not optimized for class separability. An alternative is the <u>linear discriminant analysis</u>, which does take this into account. PCA optimally minimizes reconstruction error under the L_2 norm.

PCA by COVARIANCE Method

We need to find a dxd orthonormal transformation matrix W^T, such that:

with the constraint that: Cov(Y) is a diagonal matrix, and $W^{-1} = W^{T}$.

$$Y = W^T X$$

$$COV(Y) = E[YY^{T}] = E[(W^{T}X)(W^{T}X)^{T}]$$
$$= E[(W^{T}X)(X^{T}W)] = W^{T}E[XX^{T}]W$$
$$= W^{T}COV(X)W = W^{T}(WDW^{T})W = D$$

$$WCOV(Y) = WW^T COV(X)W = COV(X)W$$

Can you derive from the above, that:

$$[\lambda_1 W_1, \lambda_2 W_2, \dots, \lambda_d W_d] =$$

 $[COV(X)W_1, COV(X)W_2, \dots, COV(X)W_d]$

Example of PCA

Samples:
$$x_1 = \begin{bmatrix} -1 \\ 1 \\ 2 \end{bmatrix}; x_2 = \begin{bmatrix} -2 \\ 3 \\ 1 \end{bmatrix}; x_3 = \begin{bmatrix} 4 \\ 0 \\ 3 \end{bmatrix}; X_3 = \begin{bmatrix} -1 & -2 & 4 \\ 1 & 3 & 0 \\ 2 & 1 & 3 \end{bmatrix}$$

3-D problem, with N = 3.

Each column is an observation (sample) and each row a variable (dimension),

Mean of the samples: $\mu_{x} = \begin{bmatrix} \frac{1}{3} \\ \frac{4}{3} \\ 2 \end{bmatrix}; \quad \tilde{x}_{1} = \begin{bmatrix} -\frac{4}{3} \\ -\frac{1}{3} \\ 0 \end{bmatrix}; \quad \tilde{x}_{2} = \begin{bmatrix} -\frac{7}{3} \\ \frac{5}{3} \\ -\frac{1}{3} \\ -1 \end{bmatrix}; \quad \tilde{x}_{3} = \begin{bmatrix} \frac{11}{3} \\ -\frac{4}{3} \\ -\frac{4}{3} \\ 1 \end{bmatrix};$

<u>Method – 1</u> (easiest)

$$\tilde{X} = \begin{bmatrix} -\frac{4}{3} & -\frac{7}{3} & \frac{11}{3} \\ -\frac{1}{3} & \frac{5}{3} & -\frac{4}{3} \\ 0 & -1 & 1 \end{bmatrix}; \quad \tilde{(X \ X} \ \tilde{X} \ \tilde{)}/2 = (1/2) \begin{bmatrix} 62/2 & -\frac{25}{3} & 6 \\ -\frac{3}{3} & \frac{-25}{3} & 6 \\ -\frac{25}{3} & \frac{14}{3} & -3 \\ 6 & -3 & 2 \end{bmatrix}$$

Method – 2 (PCA defn.)

$$S_T = (\frac{1}{N-1}) \sum_{k=1}^{N} (x_k - \mu) (x_k - \mu)^T$$

C1 =		
1.7778	0.4444	0
0.4444	0.1111	0
0	0	0

$$\tilde{x}_{1} = \begin{bmatrix} -\frac{4}{3} \\ -\frac{1}{3} \\ -\frac{1}{3} \\ 0 \end{bmatrix}; \tilde{x}_{2} = \begin{bmatrix} -\frac{7}{3} \\ \frac{5}{3} \\ -1 \end{bmatrix}; \tilde{x}_{3} = \begin{bmatrix} 1\frac{1}{3} \\ -\frac{4}{3} \\ -\frac{1}{3} \\ 1 \end{bmatrix};$$

3.6667

1.0000

C2 = 5.4444 -3.8889 2.3333 -3.8889 2.7778 -1.6667 2.3333 -1.6667 1.0000

C3 = 13.4444 -4.8889 -4.8889 1.7778 -1.3333 3.6667 -1.3333

SigmaC = 20.6667 -8.3333 6.0000 -8.3333 4.6667 -3.0000 6.0000 -3.0000 2.0000

COVAR = SigmaC/2 =

10.3333	-4.1667	3.0000
-4.1667	2.3333	-1.5000
3.0000	-1.5000	1.0000

Next do SVD, to get vectors.

SVD – the theorem

Suppose M is an m-by-n matrix whose entries come from the field K, which is either the field of real numbers or the field of complex numbers. Then there exists a factorization of the form

$M = U\Sigma V^*$

where U is an <u>m-by-m</u> unitary matrix over K, the matrix Σ is <u>m-by-n</u> with nonnegative numbers on the diagonal and zeros off the diagonal, and V* denotes the conjugate transpose of V, an n<u>-by-n</u> unitary matrix over K. Such a factorization is called a <u>singular-value decomposition of M.</u>

The matrix V thus contains a set of orthonormal "input" or "analysing" basis vector directions for M.

The matrix U contains a set of orthonormal "output" basis vector directions for M. The matrix Σ contains the singular values, which can be thought of as scalar "gain controls" by which each corresponding input is multiplied to give a corresponding output.

A common convention is to order the values $\Sigma_{i,i}$ in non-increasing fashion. In this case, the diagonal matrix Σ is uniquely determined by M (though the matrices U and V are not).

For p = min(m,n) —

<u>U is m-by-p, Σ is p-by-p, and V is n-by-p</u>.

For a face image with N samples and dimension d (=w*h, very large), we have:

The array X or Xavg of size d*N (N vertical samples stacked horizontally)

Thus XX^T will be of d*d, which will be very large. To perform eigen-analysis on such large dimension is time consuming and may be erroneous.

Thus often $X^T X$ of dimension N*N is considered for eigen-analysis. Will it result in the same, after SVD? Lets check:

$$S = \tilde{X} \tilde{X}^{T} = (1/2) \begin{bmatrix} \frac{62}{3} & -\frac{25}{3} & 6\\ -\frac{25}{3} & \frac{14}{3} & -3\\ 6 & -3 & 2 \end{bmatrix} = \begin{bmatrix} 10.3333 & -4.1667 & 3.0000\\ -4.1667 & 2.3333 & -1.5000\\ 3.0000 & -1.5000 & 1.0000 \end{bmatrix}$$

$$S^m = X^T \tilde{X} =$$

 $X^m = X^T X^{-2.2222}$

0.9444

1.2778

4.6111

-5.8889

-2.2222

-5.8889

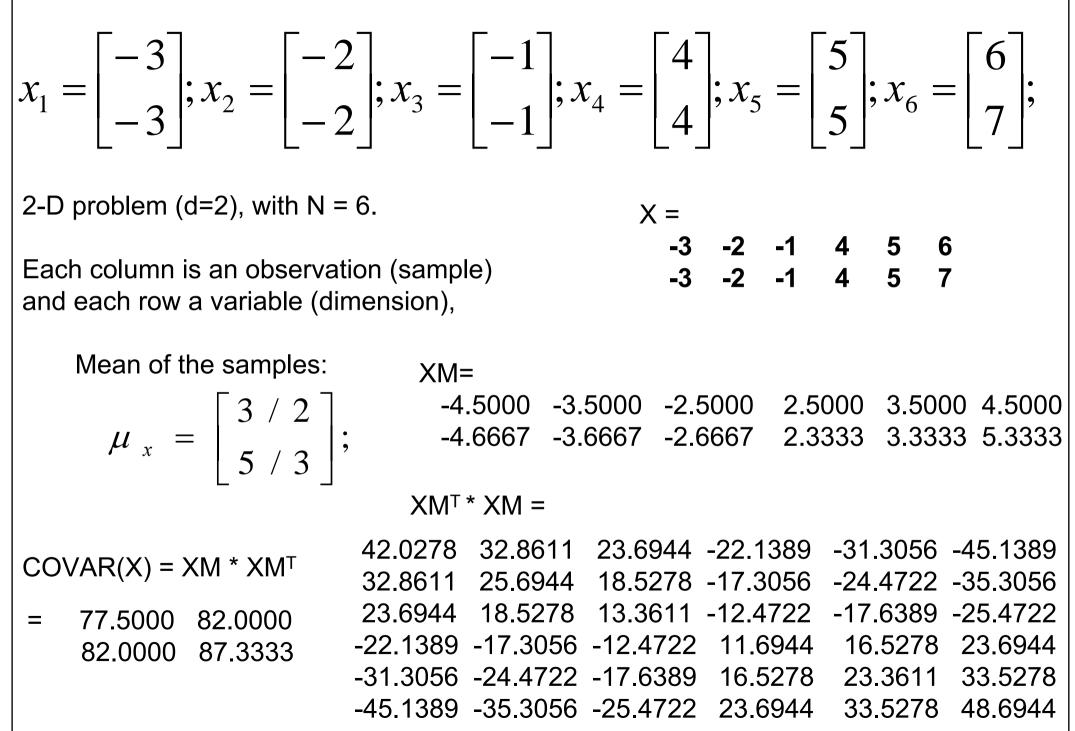
8.1111

Lets do SVD of both:

$S = X \tilde{X}^{T} =$ 10.3333 -4.1667 3.0000 -4.1667 2.3333 -1.5000 3.0000 -1.5000 1.0000	$S^{m} = X^{T} \tilde{X} =$ 0.9444 1.2778 -2.2222 1.2778 4.6111 -5.8889 -2.2222 -5.8889 8.1111
U =	U =
-0.8846 -0.4554 -0.1010 0.3818 -0.8313 0.4041 -0.2680 0.3189 0.9091	-0.2060 0.7901 0.5774 -0.5812 -0.5735 0.5774 0.7872 -0.2166 0.5774
S =	S =
$\begin{array}{cccccccccccccccccccccccccccccccccccc$	$\begin{array}{cccccccccccccccccccccccccccccccccccc$
V =	V =
-0.8846 -0.4554 0.1010 0.3818 -0.8313 -0.4041 -0.2680 0.3189 -0.9091	-0.2060 0.7901 0.5774 -0.5812 -0.5735 0.5774 0.7872 -0.2166 0.5774

Samples:

Example, where d <> N:



	XM ⁺ * XM =		
	42.0278 32.8611 23.6944 -22.1389 -31.3056 -45.1389		
COVAR(X) = XM * XM [⊤]	32.8611 25.6944 18.5278 -17.3056 -24.4722 -35.3056		
77 5000 00 0000	23.6944 18.5278 13.3611 -12.4722 -17.6389 -25.4722		
= 77.5000 82.0000	-22.1389 -17.3056 -12.4722 11.6944 16.5278 23.6944		
82.0000 87.3333	-31.3056 -24.4722 -17.6389 16.5278 23.3611 33.5278		
	-45.1389 -35.3056 -25.4722 23.6944 33.5278 48.6944		
U =			
0 6956 0 7290	U =		
-0.6856 -0.7280 -0.7280 0.6856	-0.5053 -0.1469 -0.7547 0.3882 0.0214 0.0486		
-0.7200 0.0000	-0.3951 -0.0654 0.3632 0.0984 -0.4091 0.7284		
	-0.2849 0.0162 -0.0433 -0.3456 -0.7396 -0.5002		
S =	0.2660 0.4241 -0.5083 -0.5306 -0.1150 0.4429		
	0.3762 0.5057 -0.0258 0.6601 -0.4043 -0.0539		
164.5639 0	0.5432 -0.7337 -0.1938 0.0541 -0.3293 0.1332		
0 0.2694	S =		
	164.5639 0 0 0 0 0		
	0 0.2694 0 0 0 0		
V =	0 0.0 0 0		
	0 0 0 0.0 0		
-0.6856 -0.7280	0 0 0 0 0 0 0		
-0.7280 0.6856	0 0 0 0 0.0		
	<u>V = U ??</u>		

Scatter Matrices and Separability criteria

Scatter matrices used to formulate criteria of class separability:

Within-class scatter Matrix: It shows the scatter of samples around their respective class expected

vectors.

$$S_{W} = \sum_{i=1}^{c} \sum_{x_{k} \in X_{i}} (x_{k} - \mu_{i}) (x_{k} - \mu_{i})^{T}$$

Between-class scatter Matrix: It is the scatter of the expected vectors around the mixture mean....µ is the mixture mean...

$$S_{B} = \sum_{i=1}^{c} N_{i} (\mu_{i} - \mu) (\mu_{i} - \mu)^{T}$$

Scatter Matrices and Separability criteria

Mixture scatter matrix: It is the covariance matrix of all samples regardless of their class assignments.

$$S_T = \sum_{k=1}^{N} (x_k - \mu)(x_k - \mu)^T = S_W + S_B$$

• The criteria formulation for class separability needs to convert these matrices into a number.

• This number should be larger when between-class scatter is larger or the within-class scatter is smaller.

Several Criterias are..

$$J_1 = tr(S_2^{-1}S_1)$$

$$J_3 = tr(S_1) - \mu(trS_2 - c)$$

$$J_2 = \ln \left| S_2^{-1} S_1 \right| = \ln \left| S_1 \right| - \ln \left| S_2 \right|$$

$$J_4 = \frac{trS_1}{trS_2}$$

Linear Discriminant Analysis

• Learning set is labeled – make use of this – supervised learning

• Class specific method in the sense that it tries to 'shape' the scatter in order to make it more reliable for classification.

• Select W to maximize the ratio of the between-class scatter and the within-class scatter.

Between-class scatter matrix is defined by-

$$S_{B} = \sum_{i=1}^{c} N_{i} (\mu_{i} - \mu) (\mu_{i} - \mu)^{T}$$

 μ_i mean of class X_i

 N_i is the no. of samples in class X_{i} .

Within-class scatter matrix is:

$$S_{W} = \sum_{i=1}^{c} \sum_{x_{k} \in X_{i}} (x_{k} - \mu_{i}) (x_{k} - \mu_{i})^{T}$$

Linear Discriminant Analysis

• If S_W is nonsingular W_{opt} is chosen to satisfy

$$W_{opt} = \arg \max \frac{|W^T S_B W|}{|W^T S_W W|}$$

 $W_{opt} = [w_1, w_2, \dots, w_m]$

 $\{w_i \mid i = 1, 2, ..., m\}$ is the set of eigenvectors of S_B and S_W corresponding to m largest eigen values.i.e.

$$S_B w_i = \lambda_i S_W w_i$$

• There are at most (c-1) non-zero eigen values. So upper bound of m is (c-1).

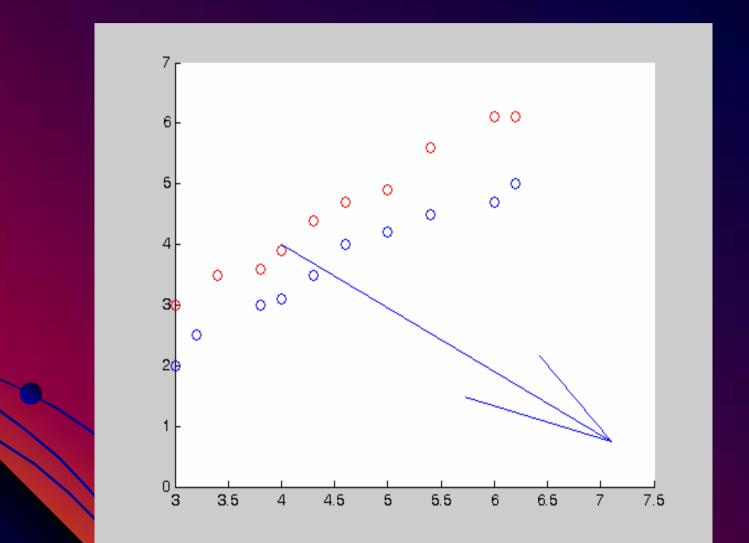
Linear Discriminant Analysis

 S_W is singular most of the time. It's rank is at most *N*-*c* Solution – Use an alternative criterion.

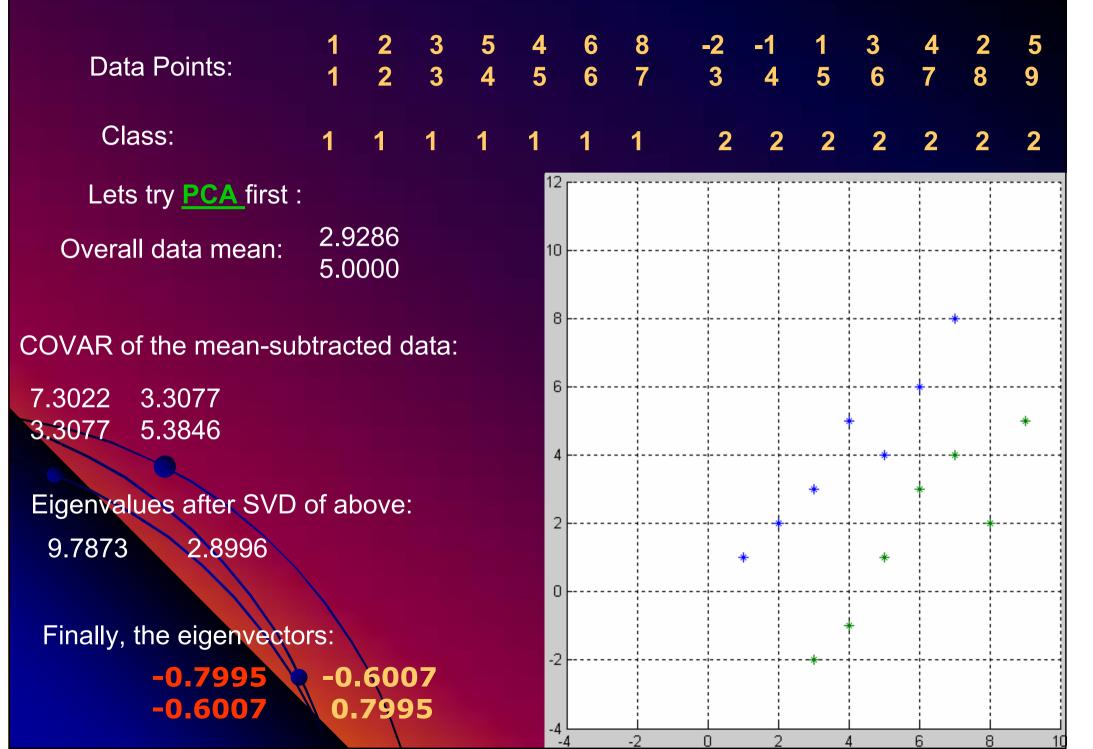
- Project the samples to a lower dimensional space.
- Use PCA to reduce dimension of the feature space to *N*-*c*.
- Then apply standard FLD to reduce dimension to *c*-1.

$$W_{opt} \text{ is given by } \qquad W_{opt} = W_{fld}^T W_{pca}^T$$
$$W_{pca} = \arg\max_{W} |W^T S_T W| \qquad W_{fld} = \arg\max_{W} \frac{|W^T W_{pca}^T S_B W_{pca} W|}{|W^T W_{pca}^T S_W W_{pca} W|}$$

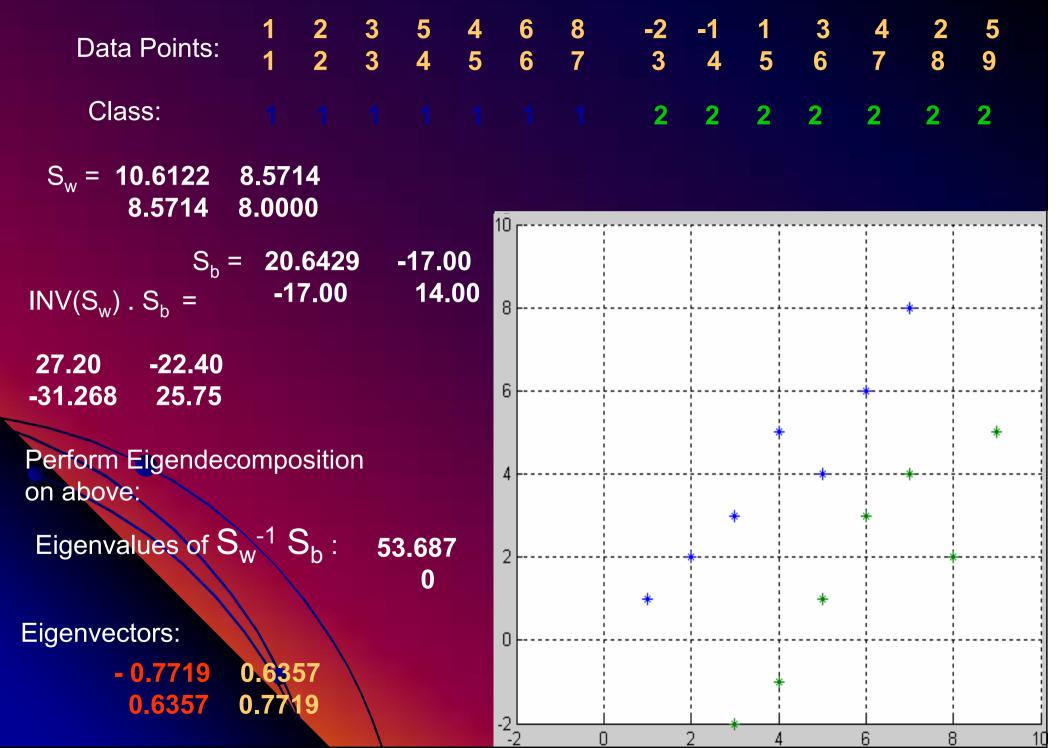
Demonstration for LDA

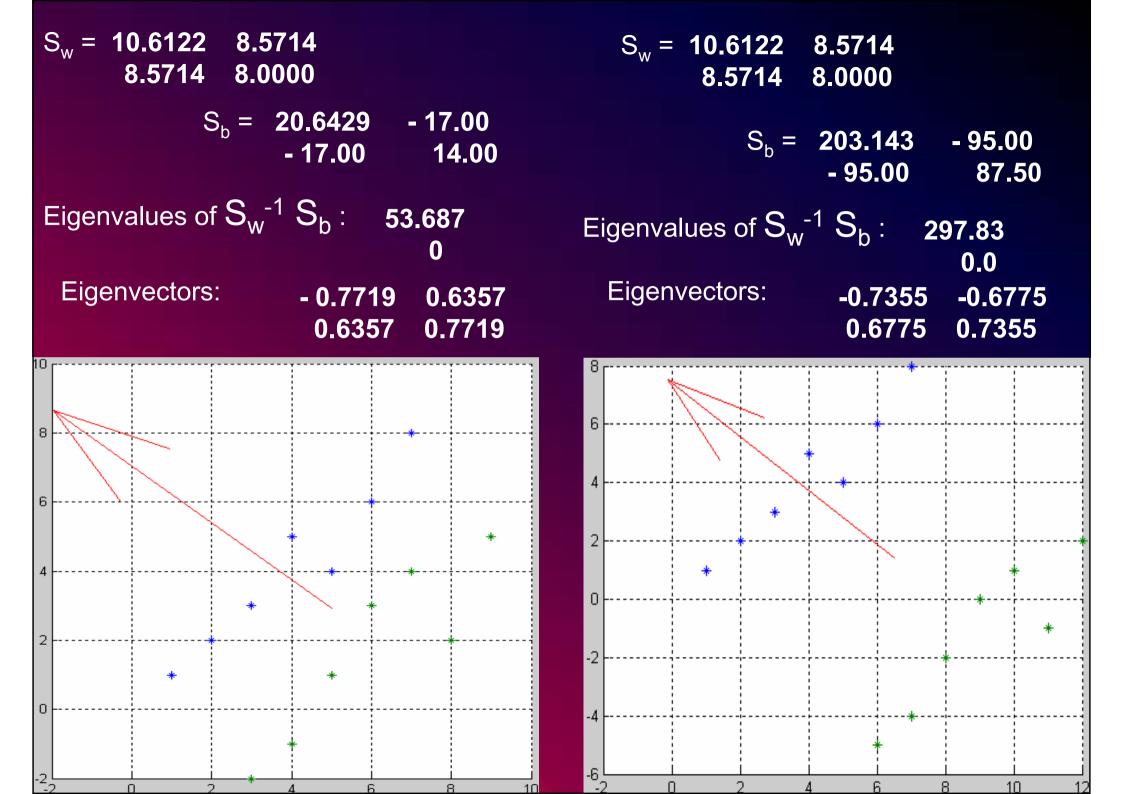


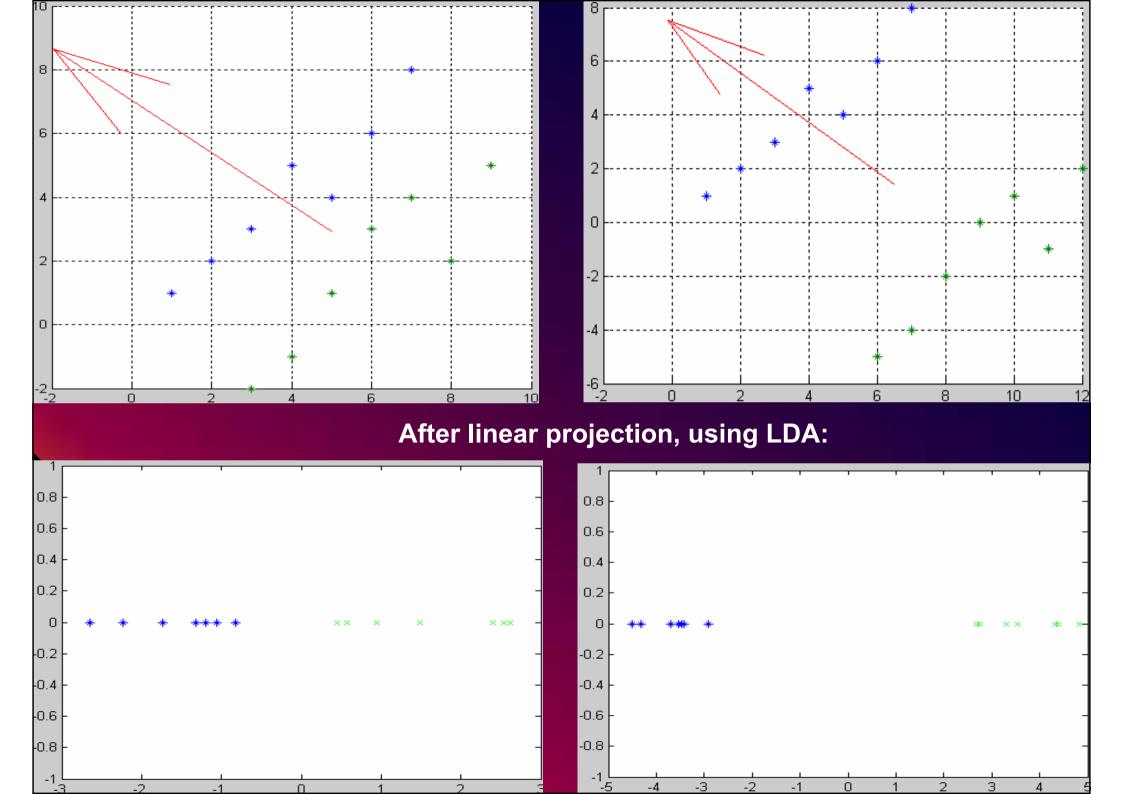
Hand workout EXAMPLE:



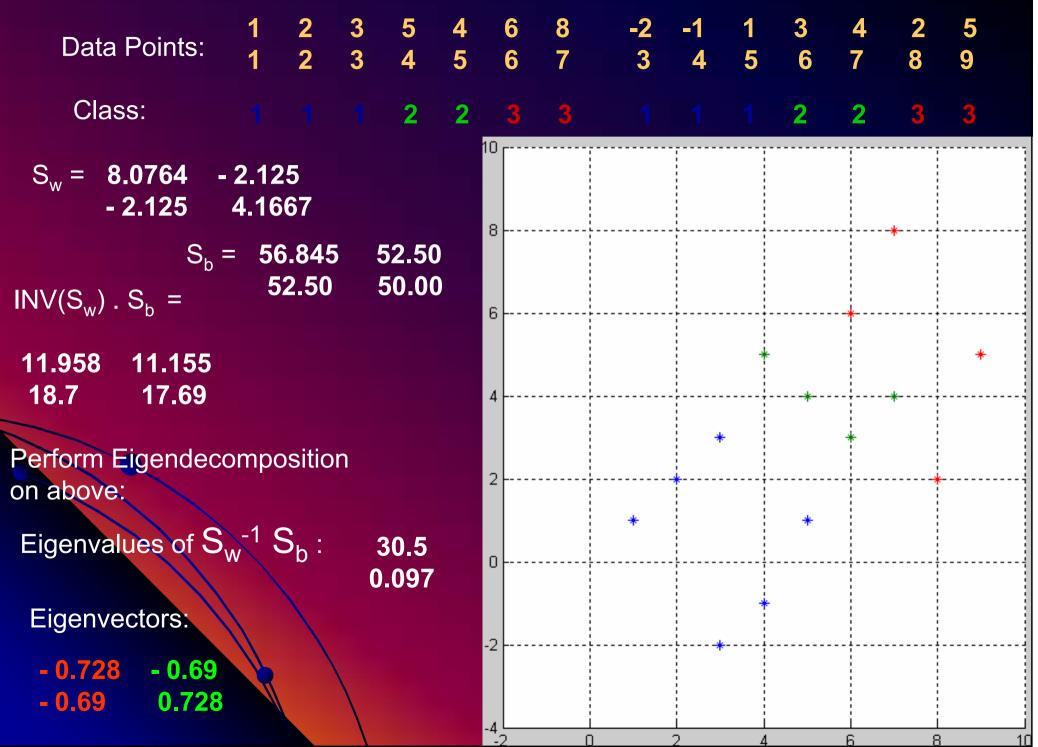
Same EXAMPLE for LDA :

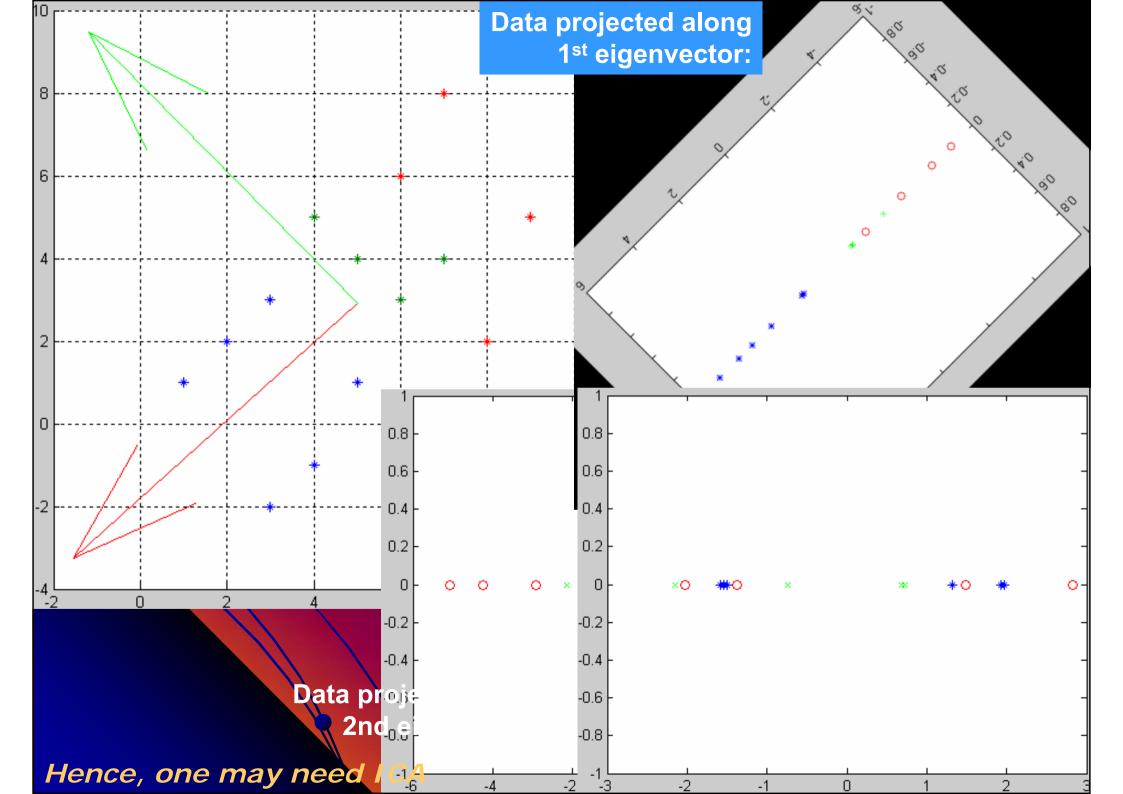






Same EXAMPLE for LDA, with C = 3:





Some of the latest advancements in Pattern recognition technology deal with:

- Neuro-fuzzy (soft computing) concepts
- Reinforcement learning
- Learning from small data sets
- Generalization capabilities
- Evolutionary Computations
- Genetic algorithms
- Pervasive computing
- Neural dynamics
- Support Vector machines

REFERENCES

 "Pattern Recognition: Statistical. Structural and Neural Approaches"; Robert J. Schallkoff; John Wiley and Sons; 1992+.

• Duda R.O., Hart P.E., D. G. Stork: Pattern Classification. John Wiley and Sons, Singapore (2001).

 Ilaria Bartolini, Paolo Ciaccia, M.I., Patella, M.: "Warp: Accurate retrieval of shapes using phase of Fourier descriptors and time warping distance". IEEE Trans. on Pattern Analysis and Machine Intelligence. Vol-27, (2005), pp. 142-147.

Statistical pattern Recognition; S. Fukunaga;