

# **Automatic Speech Recognition**

## **handout (2)**

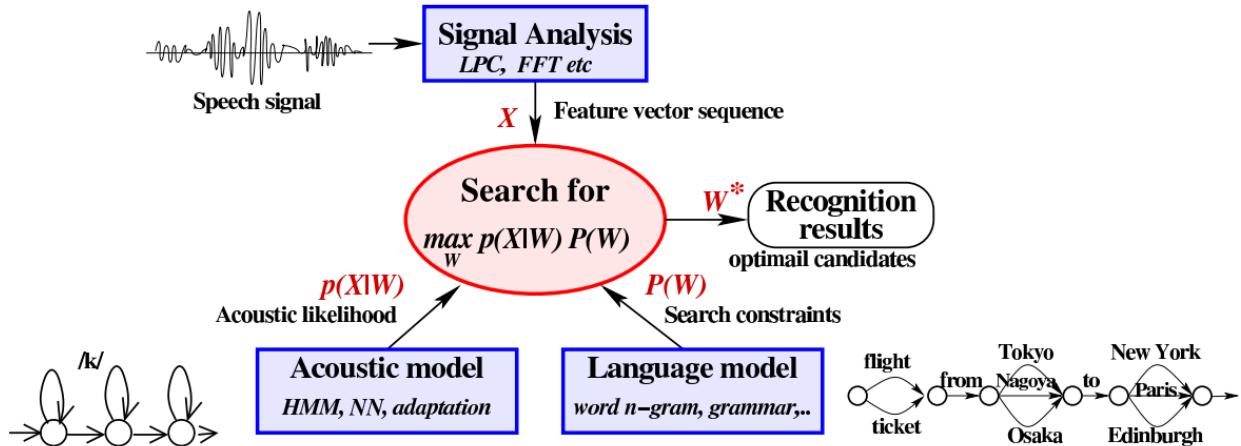
*Jan - Mar 2012*

*Revision : 1.1*

**— Statistical pattern recognition —**

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# Speech Recognition recap



(after Sagayama, "Speech Translation Telephony", 1994)

## Bayes decision rule:

$$\begin{aligned} W^* &= \arg \max_W P(W|X) = \arg \max_W \frac{p(X|W)P(W)}{p(X)} \\ &= \arg \max_W p(X|W)P(W) \end{aligned}$$

$$\begin{aligned} X &= (\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_T) : \text{feature vectors observed} \\ W &= (w_1, w_2, \dots, w_L) : \text{hypothesised words} \end{aligned}$$

# Speech Recognition recap<sub>(cont. 2)</sub>

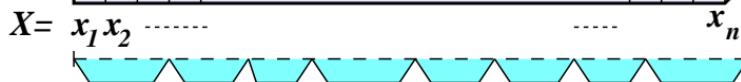
Speech signal



Spectral analysis

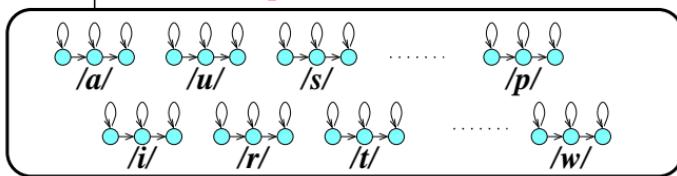


Feature vector sequence



$$p(X/\text{sayonara}) = \frac{p(X_1/\text{/s/})}{p(X_2/\text{/a/})} \frac{p(X_3/\text{/y/})}{p(X_4/\text{/o/})} \frac{p(X_5/\text{/n/})}{p(X_6/\text{/a/})} \frac{p(X_7/\text{/r/})}{p(X_8/\text{/a/})}$$

Acoustic (phone) model [HMM]



# How to calculate $p(X|/s/)$ ?

The **conditional probability** that we observe a feature sequence  $X$  for phone /s/ :

$$p(X|/s/) = p(\mathbf{x}_1, \dots, \mathbf{x}_{T_1}|/s/), \quad \mathbf{x}_i = (x_{1i}, \dots, x_{di})^t \in \mathcal{R}^d$$

We know

- HMM can be employed to calculate this. (**Viterbi algorithm, Forward / Backward algorithm**)
- HMM should be trained beforehand with a set of **training samples**. (**Baum-Welch algorithm (EM algorithm)**, Viterbi training)

To investigate these algorithms in detail, let's start with the simplest case: the length of the sequence is one ( $T = 1$ ), and the dimensionality of  $x$  is one ( $d = 1$ ).

$$p(X|/s/) \longrightarrow p(x|/s/)$$

# **How to calculate $p(x|s/)$ ?**

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$p(x|s/)$  : **conditional probability**  
**(conditional probability density function (pdf) of  $x$ )**

- We learnt that a **Gaussian / normal distribution** function can be employed to approximate the pdf by:

$$P(x|s/) = \frac{1}{\sqrt{2\pi}\sigma_{s/}} e^{-\frac{(x-\mu_{s/})^2}{2\sigma_{s/}^2}}$$

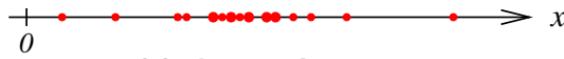
- The function has only two parameters,  $\mu_{s/}$  and  $\sigma_{s/}$
- Given a set of training samples  $\{x_1, \dots, x_N\}$ , we estimate  $\mu_{s/}$  and  $\sigma_{s/}$  by

$$\hat{\mu}_{s/} = \frac{1}{N} \sum_{i=1}^N x_i, \quad \hat{\sigma}_{s/}^2 = \frac{1}{N} \sum_{i=1}^N (x_i - \hat{\mu}_{s/})^2$$

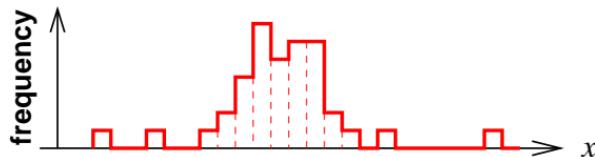
how reliable are these estimates?

how similar  $p(x|\hat{\mu}_{s/}, \hat{\sigma}_{s/}^2)$  is to the true one?

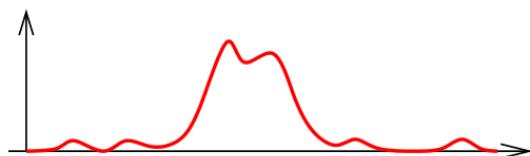
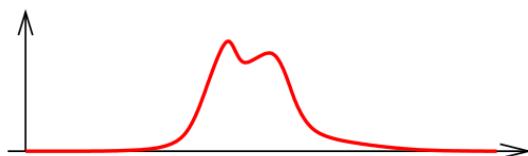
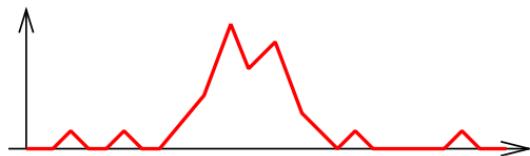
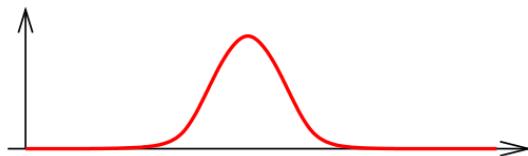
# How to estimate $p(x|s/)$ ?



(a) observation

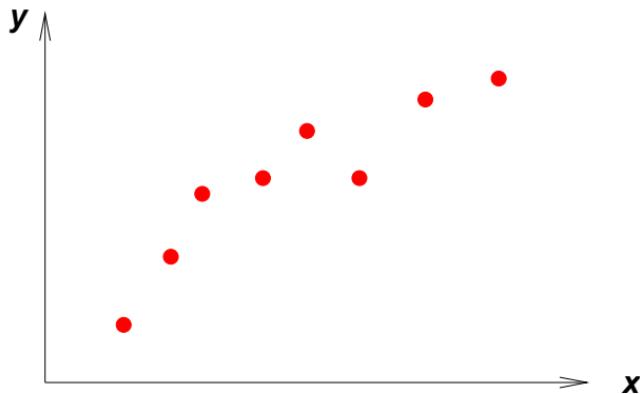


(b) histogram

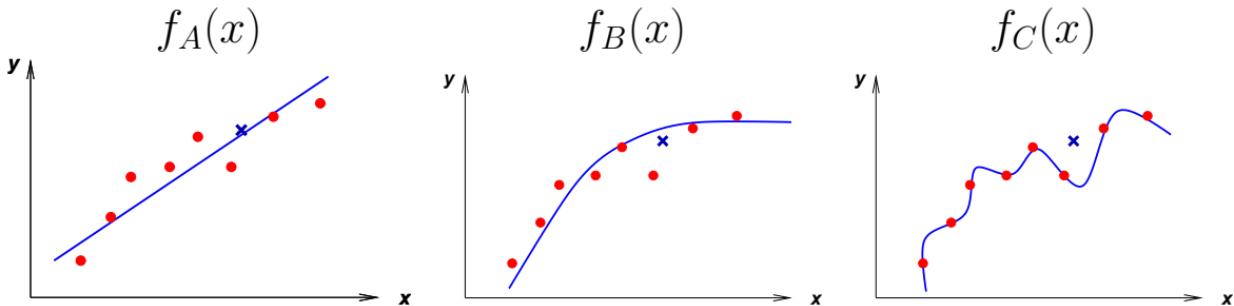


# **How to estimate $p(x|s/)$ ?**<sub>(cont. 2)</sub>

**How does  $y = f(x)$  look like?**



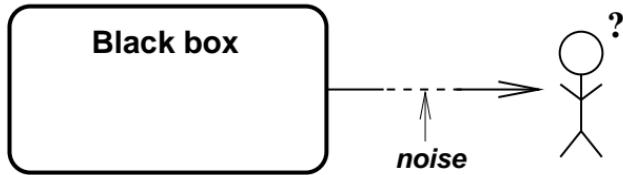
# How to estimate $p(x|s|)$ ?<sub>(cont. 3)</sub>



	$f_A(x)$	$f_B(x)$	$f_C(x)$
<b># of parameters</b>	<b>small</b>	<b>medium</b>	<b>large</b>
<b>Error on training samples</b>	<b>moderate</b>	<b>small</b>	<b>0 (zero)</b>
<b>Error on other samples</b>	<b>moderate?</b>	<b>moderate/large?</b>	<b>huge?</b>

Technical terms: **robustness, generalisation, over-fitting**

# How to estimate $p(x|s/)$ ?<sub>(cont. 4)</sub>



- What we've observed are just samples from an unobservable system. Even worse, the samples might be distorted.
- With limited amount of samples, it is hard to estimate / identify the system completely, because more than one solution can exist and we do not know which is the best one.
- To find the best one, we need more knowledge about the system.
- Estimating the system from observation is regarded as an “**inverse problem**” or “**ill-posed problem**”.
- To solve an inverse problem, we usually define an **optimisation problem** with some constraints (assumptions, models).

# *Estimating pdf*

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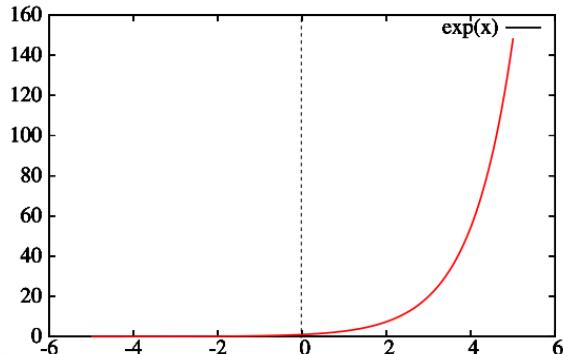
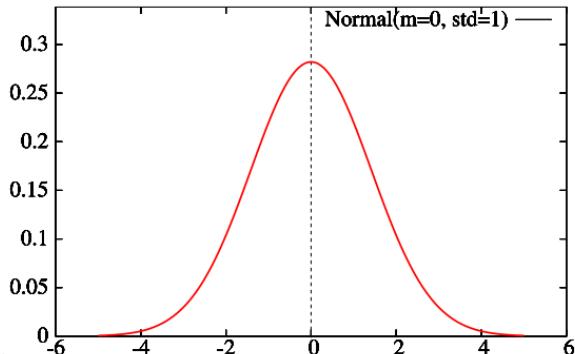
- Non-parametric approach
  - histogram
  - kernel method, Parzen window
  - neural network
- Parametric approach
  - Gaussian / normal distribution
  - Gaussian mixture model
  - etc.

# Gaussian Distribution: 1-dim case

$$p(x) = \frac{1}{\sqrt{2\pi}\sigma} e^{-\frac{(x-\mu)^2}{2\sigma^2}}$$

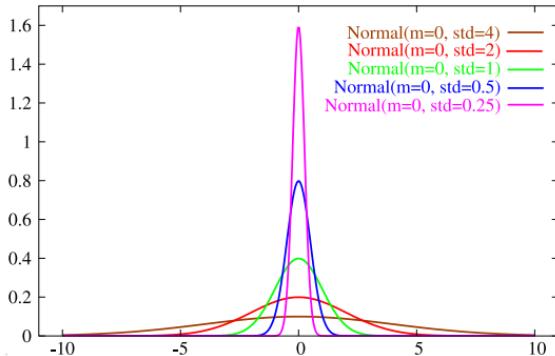
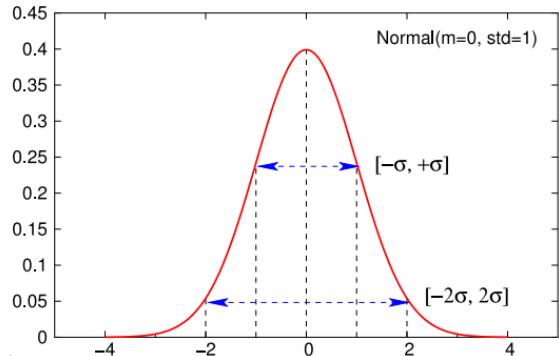
$$y = e^x$$

$$e = 2.7183\cdots$$



# Gaussian Distribution: 1-dim case

(cont. 2)



$$P(-\sigma < x < \sigma) \sim 0.683$$

$$P(-2\sigma < x < 2\sigma) \sim 0.954$$

$$P(-3\sigma < x < 3\sigma) \sim 0.997$$

$$p(x) = \frac{1}{\sqrt{2\pi}\sigma} e^{-\frac{(x-\mu)^2}{2\sigma^2}}$$

# **Gaussian Distribution: 1-dim case**<sub>(cont. 3)</sub>

$$P(x|\mu, \sigma) = \frac{1}{\sqrt{2\pi}\sigma} \exp\left(-\frac{(x-\mu)^2}{2\sigma^2}\right) \triangleq \mathcal{N}(x; \mu, \sigma)$$

**Given samples**  $\{x_1, \dots, x_N\} = \{x_n\}_{n=1}^N$ ,  
**we know how to estimate (infer) the parameters**  $\mu, \sigma$ :

$$\hat{\mu} = \frac{1}{N} \sum_{n=1}^N x_n$$

$$\hat{\sigma}^2 = \frac{1}{N} \sum_{n=1}^N (x_n - \hat{\mu})^2 = \frac{1}{N} \sum_{n=1}^N x_n^2 - \left( \frac{1}{N} \sum_{n=1}^N x_n \right)^2$$

**However,**

- Why we can say this is a right inference?
- Are there any other inferences?

# Maximum likelihood estimation

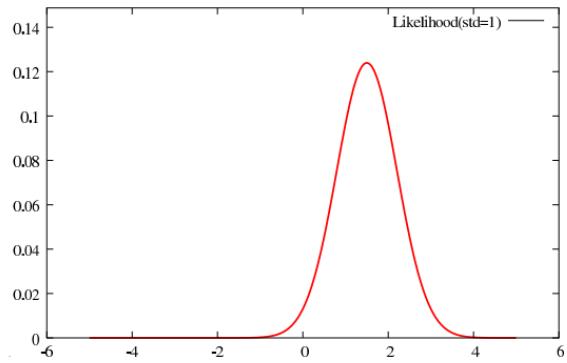
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$$p(x|\mu, \sigma) = \frac{1}{\sqrt{2\pi}\sigma} \exp\left(-\frac{(x-\mu)^2}{2\sigma^2}\right)$$

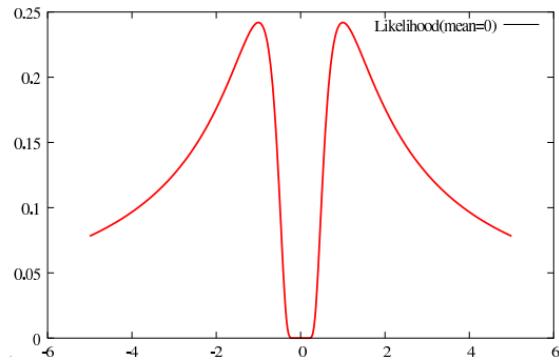
Interpretation of  $p(x|\mu, \sigma)$

- If  $\mu, \sigma$  are fixed  $\rightarrow p(x|\mu, \sigma)$  is a function of  $x$ .  
 $p(x|\mu, \sigma)$  : a probability density function.
- If  $x$  is fixed  $\rightarrow p(x|\mu, \sigma)$  is a function of  $\mu, \sigma$ .  
 $p(x|\mu, \sigma)$  : a 'likelihood' function denoted as  $L(\mu, \sigma; x)$

# Maximum likelihood estimation<sub>(cont. 2)</sub>



likelihood as a function of  $\mu$



likelihood as a function of  $\sigma$

$$L(\mu, \sigma; x) = \frac{1}{\sqrt{2\pi}\sigma} \exp\left(-\frac{(x - \mu)^2}{2\sigma^2}\right)$$

# Maximum likelihood estimation<sub>(cont. 3)</sub>

## ML estimation (MLE)

Assume  $X = \{x_n\}_{n=1}^N$  are chosen independently (i.i.d.) <sup>1</sup>.

The probability of observing  $X$  is

$$p(x_1, \dots, x_N | \mu, \sigma) = p(x_1 | \mu, \sigma) \cdots p(x_N | \mu, \sigma) = \prod_{n=1}^N p(x_n | \mu, \sigma) \\ \triangleq L(\mu, \sigma | X)$$

Find parameters  $\mu, \sigma$  that maximise the likelihood:

$$\max_{\mu, \sigma} L(\mu, \sigma | X)$$

<sup>1</sup>i.i.d: independent identically-distributed, i.e. data points that are drawn independently from the same distribution.

# **Maximum likelihood estimation** (cont. 4)

**(Solution)** Take the derivative of  $\log L(\mu, \sigma | X)$  and set it to zero.

$$\begin{aligned}\ell(\mu, \sigma | X) &= \log L(\mu, \sigma | X) = \log \prod_{n=1}^N p(x_n | \mu, \sigma) = \sum_{n=1}^N \log p(x_n | \mu, \sigma) \\ &= -N \log(\sqrt{2\pi}\sigma) - \sum_{n=1}^N \frac{(x_n - \mu)^2}{2\sigma^2}\end{aligned}$$

$$\frac{\partial \ell}{\partial \mu} = -2 \sum_{n=1}^N \frac{x_n - \mu}{2\sigma^2} = 0 \quad \Rightarrow \quad \mu = \frac{1}{N} \sum_{n=1}^N x_n$$

$$\begin{aligned}\frac{\partial \ell}{\partial \sigma} &= -N \frac{1}{\sqrt{2\pi}\sigma} + \sum_{n=1}^N \frac{(x_n - \mu)2}{\sigma^3} = 0 \\ &\Rightarrow \quad \sigma^2 = \frac{1}{N} \sum_{n=1}^N (x_n - \mu)^2\end{aligned}$$

# *Limitation of a Gaussian pdf*

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- Gaussian distribution is assumed
- does not fit more complicated distributions (bimodal, etc.)

## Solutions

- Transform features  
(coordinate **non-linear** transformation)
- Use other distribution functions
  - Poisson dist.
  - exponential dist.
  - gamma dist.
  - log-normal dist.
- Use a mixture of Gaussian distribution functions
- Use neural networks, kernel approaches

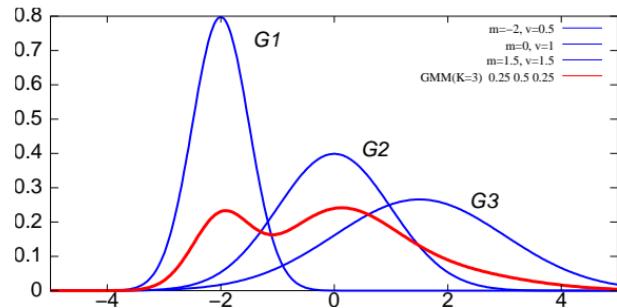
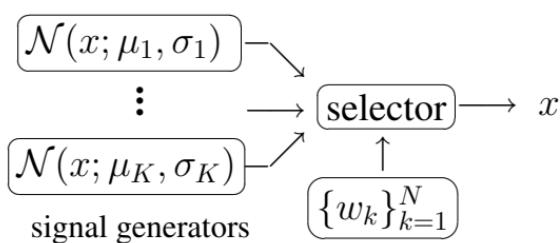
# Gaussian Mixture Model (GMM)

**Idea** Approximate a “true” distribution by more than one Gaussian distribution.

$$p(x|\Lambda) = \sum_{k=1}^K w_k p(x|\mu_k, \sigma_k) = \sum_{k=1}^K w_k \mathcal{N}(x; \mu_k, \sigma_k)$$

where  $\Lambda = \{\lambda_k\}_{k=1}^K = \{\mu_k, \sigma_k\}_{k=1}^K \quad \dots$  a set of model parameters  
 $\sum_{k=1}^K w_k = 1$

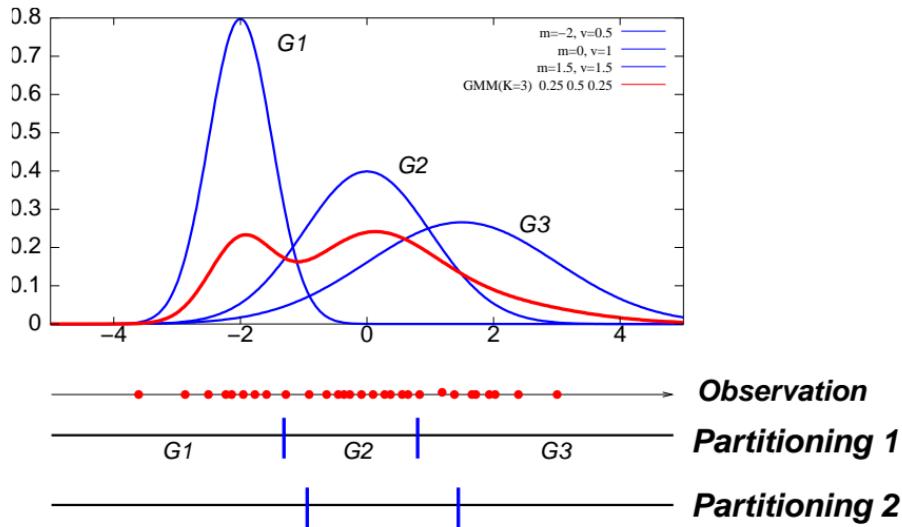
**View as a generative-model**



# Training GMM

Given a set of training samples  $x_1, x_2, \dots, x_N$ ,  
assuming a mixture of  $K$  Gaussians,  
we don't know from which Gaussian model each  $x_i$  came.

⇒ ML training can not be applied directly !



## Training GMM<sub>(cont. 2)</sub>

**Solution** Assume each  $x_n$  has a label  $y_n$  showing from which model the data came.

i.e.  $y_n$  takes an index of Gaussian model  $1, \dots, K$ .

<b>observations</b>	$x_1$	$x_2$	$\cdots$	$x_N$
<b>labels (unseen)</b>	$y_1$	$y_2$	$\cdots$	$y_N$

Then we would be able to use ML estimation recursively to find the optimal partitioning.

**Approach A** (hard  $k$ -means clustering) ... assume each  $x_n$  belongs to an exact Gaussian  $y_n$ .

**Approach B** (soft  $k$ -means clustering) ... assume each  $x_n$  belongs to Gaussian  $k$  with probability  $P(y_n = k)$ .  
( $y_n$  is treated as a random variable)

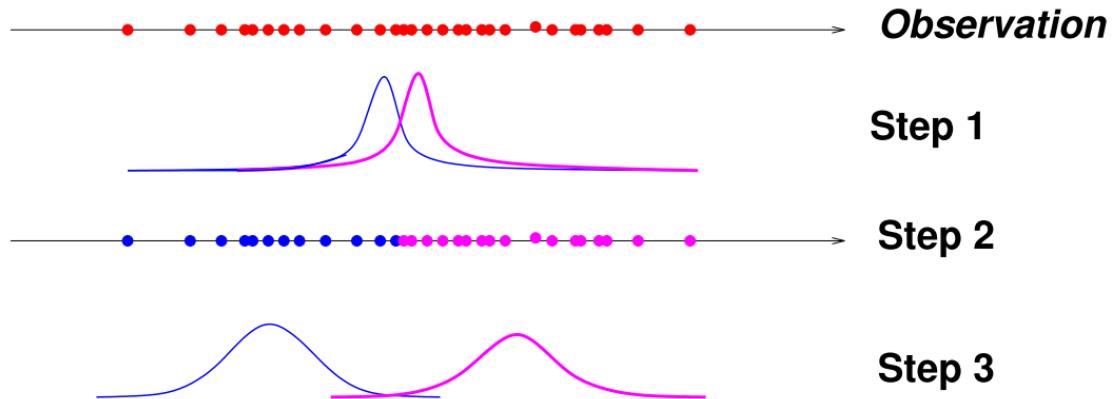
## Training GMM by $k$ -means clustering (hard version)

**Step 1:** Initialisation. Set model parameters  $\Lambda$  arbitrarily.

**Step 2:** Clustering:  $y_n = \arg \max_k p(x_n | \lambda_k)$  for  $k = 1, \dots, N$

**Step 3:** Re-estimation: estimate model parameters  $\Lambda$ .

**Step 4:** Iteration: Repeat steps 2 and 3 until a predefined convergence criterion is satisfied. (estimate  $\{w_k\}$  when exits.)



## Training GMM by $k$ -means clustering (soft version)

Problem of the hard  $k$  means clustering: non guarantee of ML.

The **EM algorithm** enables us to maximise the likelihood

$$L(\Lambda|X) = \prod_{n=1}^N p(x_n|\Lambda)$$

by iteratively maximising  $Q$  function with respect to  $\Lambda^{(t+1)}$ :

$$Q(\Lambda^{(t)}, \Lambda^{(t+1)}) = \sum_Y P(Y|X, \Lambda^{(t)}) \log p(X, Y|\Lambda^{(t+1)})$$

This is the expectation of log-likelihood from **complete data**  $(X, Y)$ .  
To maximise  $Q$ , we take the derivative of  $Q$  with respect to each parameter and set it to zero.

## **The EM algorithm** [Dempster, Laird, and Rubin, 1977]

**Step 1: Initialisation:** Set  $t = 0$ , and choose an initial estimate  $\Lambda^{(0)}$ .

**Step 2: E-Step:** Compute  $Q(\Lambda^{(t)}, \Lambda)$  based on current parameter  $\Lambda^{(t)}$ .

**Step 3: M-Step:** Compute  $\Lambda^* = \arg \max_{\Lambda} Q(\Lambda^{(t)}, \Lambda)$  to maximise  $Q$ .

**Step 4: Iteration:** Set  $\Lambda^{(t+1)} = \Lambda^*$  and  $t \leftarrow t + 1$ , repeat from Step 2 until convergence.

**(Note)**

- The algorithm is a generalisation of MLE, when we have incomplete data (i.e. some data are unobservable (hidden)).
- The algorithm gives just a general idea, and actual implementation differs depending on the problems.

# EM algorithm for GMM

$$\begin{aligned} Q(\Lambda^{(t)}, \Lambda^{(t+1)}) &= \sum_{n=1}^N Q_i(\Lambda^{(t)}, \Lambda^{(t+1)}) = \sum_{n=1}^N \sum_{y_n} P(y_n|x_n, \Lambda^{(t)}) \log p(x_n, y_n|\Lambda^{(t+1)}) \\ &= \sum_{n=1}^N \sum_{y_n} \frac{p(x_n, y_n|\Lambda^{(t)})}{p(x_n|\Lambda^{(t)})} \log p(x_n, y_n|\Lambda^{(t+1)}) \\ &= \sum_{k=1}^K \gamma_k \log w_k^{(t+1)} + \sum_{k=1}^K Q_\pi(\Lambda^{(t)}, \lambda_k^{(t+1)}) \end{aligned}$$

where

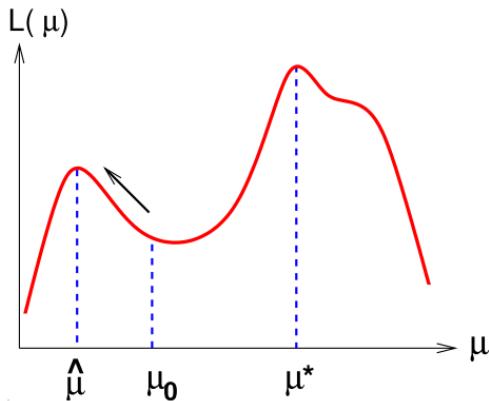
$$\begin{aligned} \gamma_k &= \sum_{n=1}^N \gamma_k^n, & \gamma_k^n &= \frac{w_k^{(t)} p(x_n|\lambda_k^{(t)})}{p(x_n|\Lambda^{(t)})} \\ Q_\pi(\Lambda^{(t)}, \lambda_k^{(t+1)}) &= \sum_{n=1}^N \gamma_k^n \log p(x_n|\lambda_k^{(t+1)}) \\ w_k^{(t+1)} &= \gamma_k / \sum_{k=1}^K \gamma_k = \gamma_k / N \\ m_k^{(t+1)} &= \sum_{n=1}^N \gamma_k^n x_n / \sum_{n=1}^N \gamma_k^n, & \sigma_k^{(t+1)} &= \sum_{n=1}^N \gamma_k^n (x - m_k^{(t)})^2 / \sum_{n=1}^N \gamma_k^n \end{aligned}$$

## *Limitation of ML-trained GMM*

- **Local optimum problem** (the EM algorithm does not guarantee the global optimum)
- **No mechanism for determining  $K$  (# of mixtures).**

# Initial value problem of GMM

- The EM iteration converges at a local maximum, which might be different from the global maximum.
- The performance of the algorithm highly relies on how appropriate initial value the algorithm starts from.
- As we increase the number of Gaussian mixture components, the chance we suffer the problem increases.

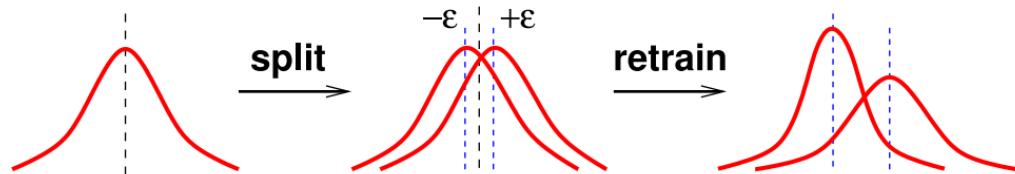


# Initial value problem of GMM<sub>(cont. 2)</sub>

Techniques to avoid the local maximum problem

**Method-A:** Use  $k$ -means clustering (hard) to find initial partitions of data, and calculate initial parameters of GMM.

**Method-B:** Increase the mixture number  $k$  of GMM gradually.  
(start from  $k = 1$ , and then  $k \leftarrow k * 2.$  )



**Fig. Mixture splitting**

# ***Remaining problems with GMM***

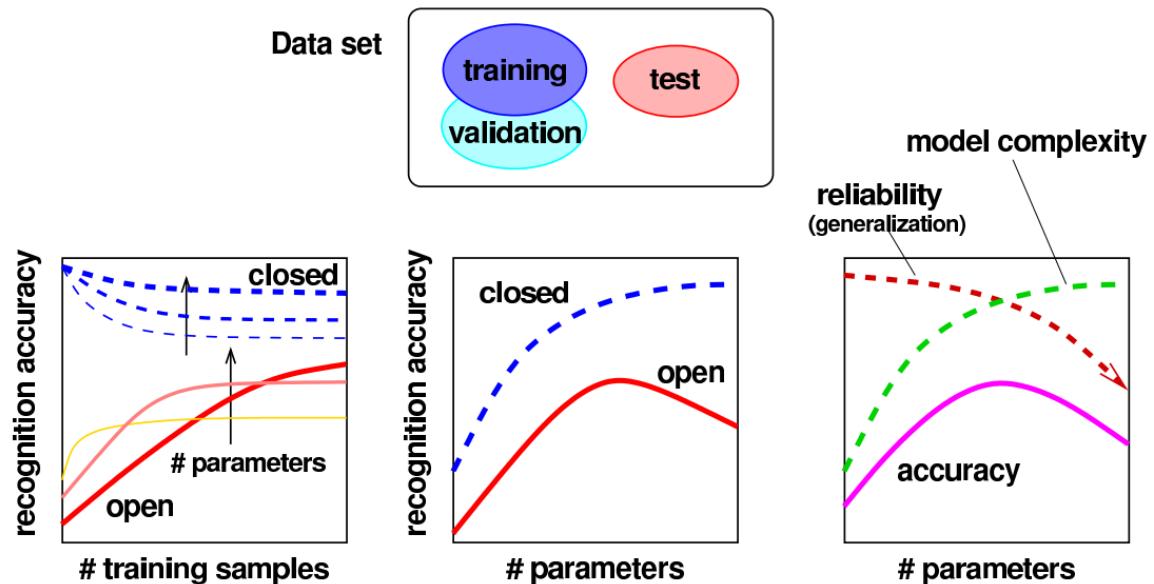
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- Avoiding the local maximum problem is not enough
- There is another problem: “**over-fitting problem**”
  - the likelihood on training data increases, as  $k$  (# of mixtures) increases.
  - The ML training does not tell us what  $k$  is optimal.

## **Solutions**

- use model selection technique such as
  - **AIC** (Akaike's information criterion)
  - **MDL** (minimum description length)  
→ “**Occam's razor**”:“**Entia non sunt multiplicanda praeter necessitatem**”  
(accepts the simplest explanation that fits the data)
- use **validation data**

# Remaining problems with GMM<sub>(cont. 2)</sub>



# **Maximum a posterior (MAP) estimation**

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$X$  : a set of training data  $\{x_n\}_{n=1}^N$

$\Lambda$  : a set of parameters to estimate

**ML estim.**  $\Lambda^* = \arg \max_{\Lambda} p(X|\Lambda)$

**MAP estim.**  $\Lambda^* = \arg \max_{\Lambda} p(\Lambda|X) = \arg \max_{\Lambda} p(X|\Lambda)p(\Lambda)$

$p(\Lambda)$  : a priori distribution over  $\Lambda$

$p(\Lambda|X)$  : a posteriors distribution after observing  $X$

$$p(\Lambda|X) = \frac{p(X|\Lambda)p(\Lambda)}{p(X)} = \frac{p(X|\Lambda)p(\Lambda)}{\int p(X|\Lambda')p(\Lambda')d\Lambda'}$$

- The object function takes the reliability of parameters, i.e. the a priori distribution, into account
- Regarded as a penalised (regularised) version of MLE
- How to define/calculate the prior  $p(\Lambda)$  ?  
⇒ several approaches have been proposed

# Bayesian estimation

p.d.f. estimation based on

- point estimation

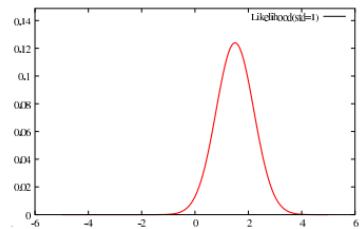
$$\hat{p}(\mathbf{x}|X) = p(\mathbf{x}|X, \Lambda^*), \quad \Lambda^* = \begin{cases} \max_{\Lambda} p(X|\Lambda), & \text{MLE} \\ \max_{\Lambda} p(X|\Lambda)p(\Lambda), & \text{MAP} \end{cases}$$

assuming parameter  $\Lambda$  is fixed but unknown.

- interval estimation (Bayesian estimation)

$$\hat{p}(\mathbf{x}|X) = \int p(\mathbf{x}|X, \Lambda)p(\Lambda|X)d\Lambda$$

$$\text{where } p(\Lambda|X) = \frac{p(X|\Lambda)p(\Lambda)}{\int p(X|\Lambda)p(\Lambda)d\Lambda}$$



assuming the parameters  $\Lambda$  are random variables with a prior distribution  $p(\Lambda)$

# Multivariate Gaussian pdf

## One variable

**data:**  $x_n = (x_n)$

$$\mu = \frac{1}{N} \sum_{i=1}^N x_i, \quad \sigma^2 = \frac{1}{N} \sum_{i=1}^N (x_i - \mu)^2$$

## Two variables

$$\mathbf{x}_n = (x_{1n}, x_{2n})^t = \begin{pmatrix} x_{1n} \\ x_{2n} \end{pmatrix}, \quad \boldsymbol{\mu} = \begin{pmatrix} \mu_1 \\ \mu_2 \end{pmatrix}, \quad \Sigma = \begin{pmatrix} \sigma_{11} & \sigma_{12} \\ \sigma_{21} & \sigma_{22} \end{pmatrix}$$

$$\mu_1 = \frac{1}{N} \sum_{i=1}^N x_{1n}, \quad \mu_2 = \frac{1}{N} \sum_{i=1}^N x_{2n}$$

$$\sigma_{11} = \frac{1}{N} \sum_{i=1}^N (x_{1i} - \mu_1)^2, \quad \sigma_{22} = \frac{1}{N} \sum_{i=1}^N (x_{2i} - \mu_2)^2$$

$$\sigma_{12} = \frac{1}{N} \sum_{i=1}^N (x_{1i} - \mu_1)(x_{2i} - \mu_2) = \sigma_{21}$$

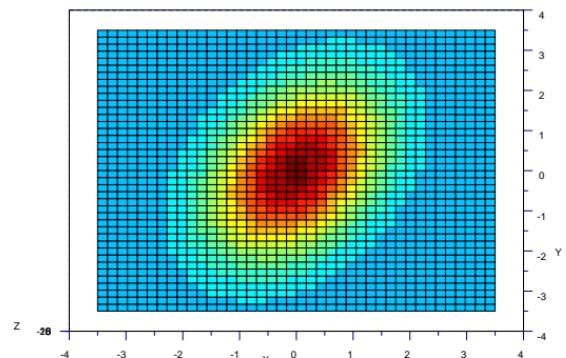
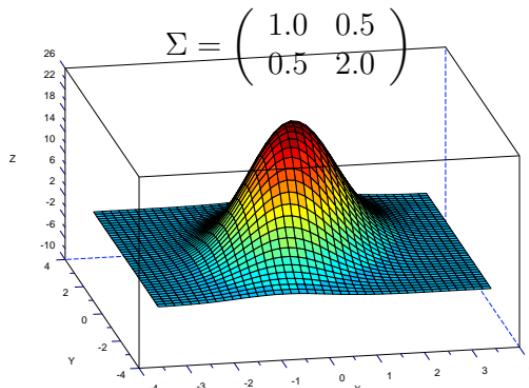
# Multivariate Gaussian pdf<sub>(cont. 2)</sub>

**mean vector:**  $\boldsymbol{\mu} = (\mu_1, \mu_2)^t$

**covariance matrix:**  $\Sigma = \begin{pmatrix} \sigma_{11} & \sigma_{12} \\ \sigma_{21} & \sigma_{22} \end{pmatrix}$

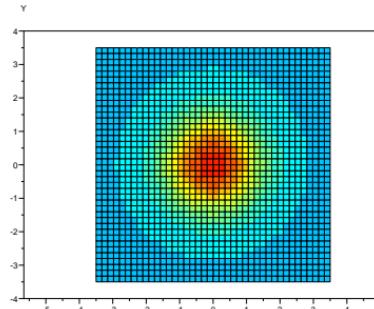
$$p(\mathbf{x}|\boldsymbol{\mu}, \Sigma) = \frac{1}{(2\pi)^{d/2} |\Sigma|^{1/2}} \exp \left( -\frac{1}{2} (\mathbf{x} - \boldsymbol{\mu})^t \Sigma^{-1} (\mathbf{x} - \boldsymbol{\mu}) \right)$$

where  $d$  denotes dimensions (here  $d = 2$ ).

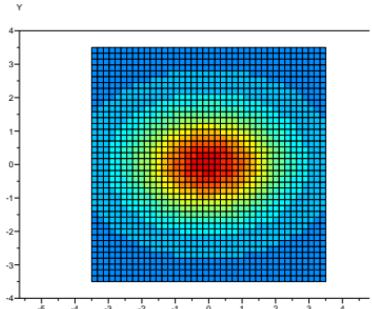


# Multivariate Gaussian pdf<sub>(cont. 3)</sub>

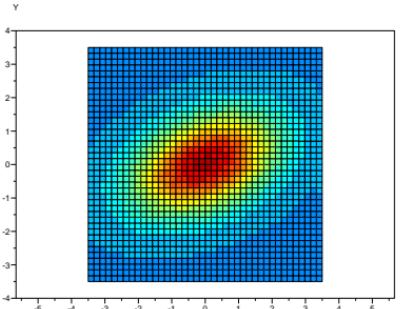
$$\Sigma = \begin{pmatrix} 1.0 & 0.0 \\ 0.0 & 1.0 \end{pmatrix}$$



$$\Sigma = \begin{pmatrix} 2.0 & 0.0 \\ 0.0 & 1.0 \end{pmatrix}$$



$$\Sigma = \begin{pmatrix} 2.0 & 0.5 \\ 0.5 & 1.0 \end{pmatrix}$$

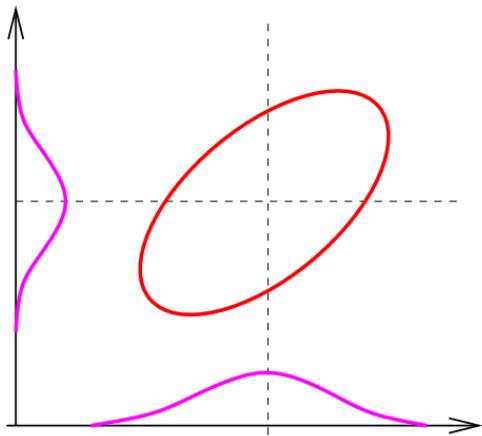


$x_1$  and  $x_2$  are called **uncorrelated** if  $\sigma_{12} = 0$ .

# Gaussian pdf with diagonal covariances

Gaussian pdf  
with a full-covariance

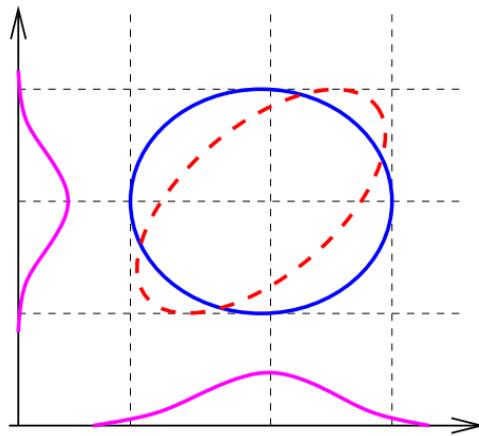
$$\Sigma = \begin{pmatrix} \sigma_{11} & \sigma_{12} \\ \sigma_{21} & \sigma_{22} \end{pmatrix}$$



Gaussian pdf  
with diagonal-covariances

$$\Sigma = \begin{pmatrix} \sigma_{11} & 0 \\ 0 & \sigma_{22} \end{pmatrix}$$

$$p(\mathbf{x}|\boldsymbol{\mu}, \Sigma) = p(x_1|\mu_1, \sigma_{11}) p(x_2|\mu_2, \sigma_{22})$$

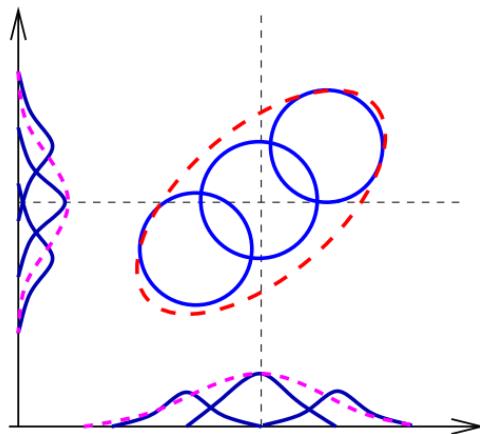


GMM with diag. covs. can model dependencies (correlations) !

Gaussian pdf  
with full-covariance



GMM pdf  
with diagonal-covariance



## Gaussian pdf with diagonal covariances<sub>(cont. 3)</sub>

**Question** why we normally use GMM with diagonal covariances rather than the one with full covariances?

# What has not been discussed yet

$$W^* = \arg \max_W P(W|X) = \arg \max_W p(X|W)P(W)$$

$X = (\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_T)$  : feature vectors observed  
 $W = (w_1, w_2, \dots, w_L)$  : hypothesised words

## ■ Assumptions so far

- $T = 1$  :  $X$  consists of a single frame
  - $L = 1$  and  $W = /s/$  (single phoneme  $/s/$ )
- ## ■ Calculation of $P(X|W)$ and $P(W)$ under those conditions:
- $X$  is a sequence of feature vectors
  - $W$  is a sequence of words/phonemes ( $L$  is unknown)
  - boundaries of words/phonemes (+ silence) are unknown

# Basic idea of estimating error rate

## 1. Prepare two subsets of the dataset:

**training set** to train the recogniser

**test(validation) set** to evaluate recognition performance

\* these two subsets should be mutually exclusive (no overlap between them)

## 2. Train the recogniser using the training set.

## 3. Carry out a recognition experiment of the recogniser on the test set, and calculate the recognition error:

$$R_{\text{emp}} = \frac{\text{\# of misrecognised data}}{\text{\# of input data}} \quad \dots \quad \left( \begin{array}{l} \text{test set error, or} \\ \text{empirical error} \end{array} \right)$$

## 4. Repeat the recognition experiment with changing training and test sets to estimate the expectation of error:

**generalisation error:**  $R = E[R_{\text{emp}}]$

[http://en.wikipedia.org/wiki/Generalization\\_error](http://en.wikipedia.org/wiki/Generalization_error)

# *k-fold Cross-Validation (CV)*

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**Step 1** split the dataset randomly into  $k$  disjoint sets of equal size.

**Step 2** train and test the recogniser  $k$  times: each time with choosing one set for validation and the remaining  $k - 1$  sets for training.

**Step 3** Estimate the performance by taking the mean of the  $k$  errors.

<b>k-fold CV</b>	<b>subsets</b>
<b>2-fold CV</b>	<b>Set2-0, Set2-1</b>
<b>5-fold CV</b>	<b>Set5-0, Set5-1, Set5-2, Set5-3, Set5-4</b>

Similar techniques: leave-one-out, held-out, jackknife, bootstrap

References:

<http://en.wikipedia.org/wiki/Cross-validation>

<http://www.faqs.org/faqs/ai-faq/neural-nets/part3/section-12.html>

# References

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- Books:

- Christopher M. Bishop, “Pattern Recognition and Machine Learning”, ISBN 0387310738, Springer-Verlag, 2006.

- Web pages:

## Normal (Gaussian) Distribution

- [http://en.wikipedia.org/wiki/Normal\\_distribution](http://en.wikipedia.org/wiki/Normal_distribution)

## Maximum Likelihood Estimation (MLE)

- [http://en.wikipedia.org/wiki/Maximum\\_likelihood](http://en.wikipedia.org/wiki/Maximum_likelihood)

## GMM

- [http://en.wikipedia.org/wiki/Mixture\\_model](http://en.wikipedia.org/wiki/Mixture_model)

## EM algorithm

- [http://en.wikipedia.org/wiki/Expectation-maximization\\_algorithm](http://en.wikipedia.org/wiki/Expectation-maximization_algorithm)