

Automatic Speech Recognition handout (2)

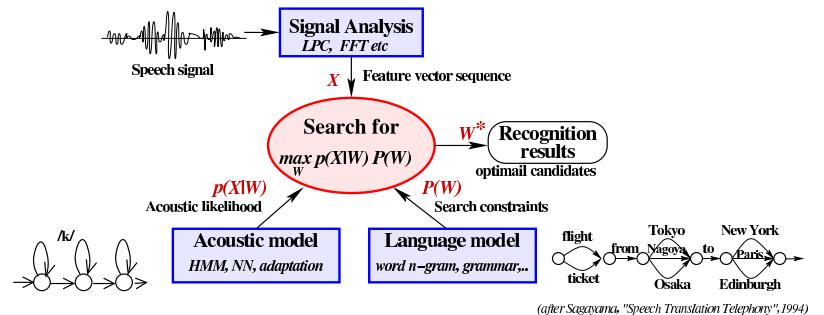
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Revision : 1.1

— Generative models and training algorithms —

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Speech Recognition recap (cont. 2)



Bayes decision rule:

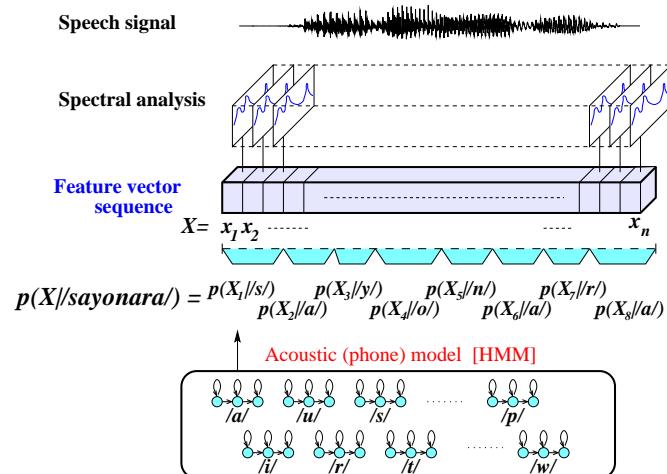
$$\begin{aligned} W^* &= \arg \max_W P(W|X) = \arg \max_W \frac{p(X|W)P(W)}{p(X)} \\ &= \arg \max_W p(X|W)P(W) \end{aligned}$$

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

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Speech Recognition recap



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 assuming the simplest case: the length of the sequence is one ($T_1 = 1$), and the dimensionality of x is one ($d = 1$).

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

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II:3

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II:1

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

$p(x|s/)$: 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, μ_s and σ_s
- Given a set of training samples $\{x_1, \dots, x_N\}$, we estimate μ_s and σ_s by

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

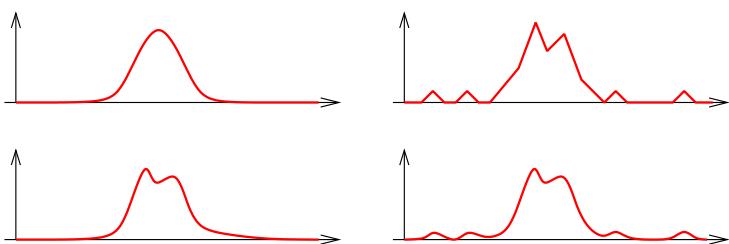
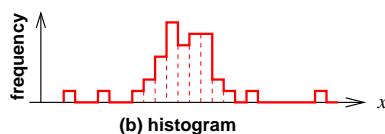
But, how reliable are these estimates ?
 how similar $p(x|\hat{\mu}_s, \hat{\sigma}_s)$ is to the true one ?

Why learning is so hard ?

Who is this person ?

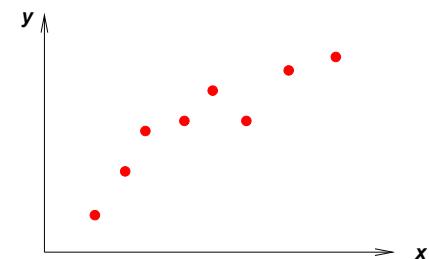


How to estimate $p(x|s/)$?

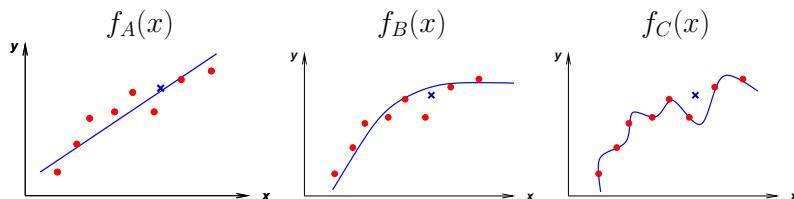


Why learning is so hard ?_(cont. 2)

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



Why learning is so hard ? (cont. 3)



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

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

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Estimating pdf

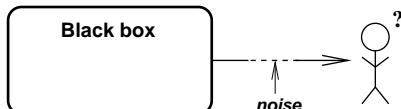
■ Non-parametric approach

- histogram
- kernel method, Parzen window
- neural network

■ Parametric approach

- Gaussian / normal distribution
- Gaussian mixture model
- etc.

Why learning is so hard ? (cont. 4)



- 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).

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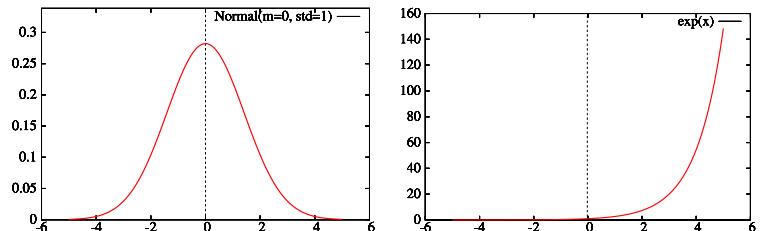
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Gaussian distribution

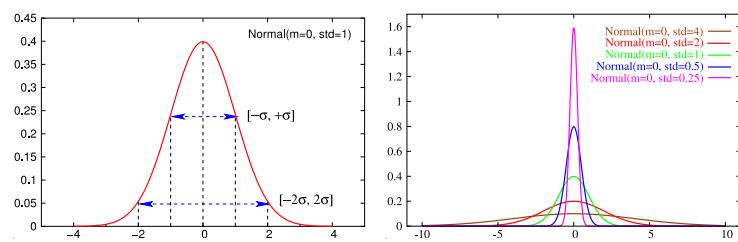
$$p(x) = \frac{1}{\sqrt{2\pi}\sigma} e^{-\frac{(x-\mu)^2}{2\sigma^2}}$$

$$y = e^x$$

$e = 2.7183\dots$



Gaussian distribution (cont. 2)



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

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

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

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Maximum likelihood estimation

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

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

- If μ, σ 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 μ, σ .
 $p(x|\mu, \sigma)$: a 'likelihood' function denoted as $L(\mu, \sigma; x)$

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Gaussian Distribution: 1-dim case

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

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

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

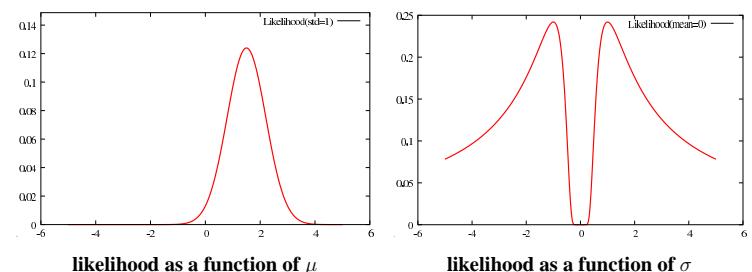
But,

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

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Maximum likelihood estimation (cont. 2)



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

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Maximum likelihood estimation (cont. 3)

ML estimation (MLE)

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

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 μ, σ that maximise the likelihood:

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

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

Limitation of a Gaussian pdf

- 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

Maximum likelihood estimation (cont. 4)

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

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

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

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

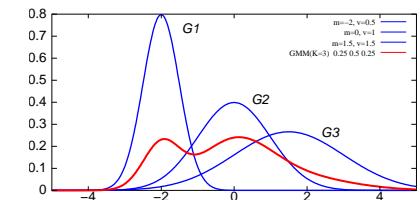
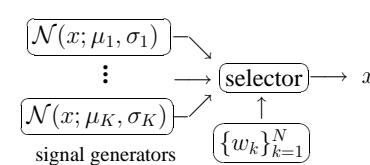
Gaussian Mixture Model (GMM)

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

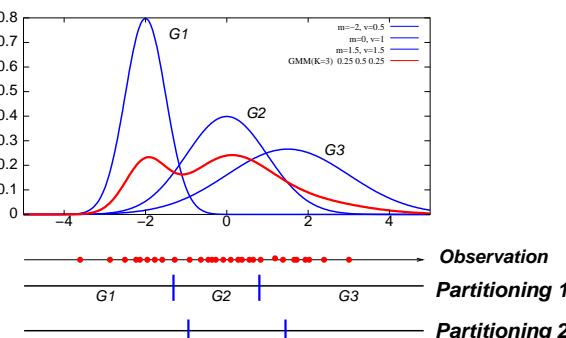
View as a generative-model



Training GMM

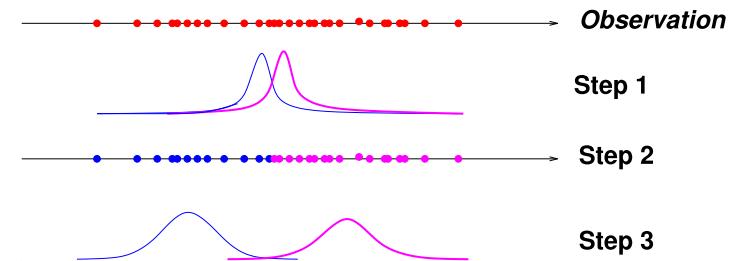
Given a set of training samples x_1, x_2, \dots, x_N ,
assuming a mixture of K Gaussians,
we don't know from which Gaussian model each x_i came.

⇒ ML training can not be applied directly !



Training GMM by k -means clustering (hard version)

- Step 1: Initialisation. Set model parameters Λ 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 Λ .
- Step 4: Iteration: Repeat steps 2 and 3 until a predefined convergence criterion is satisfied. (estimate $\{w_k\}$ when exists.)



Training GMM (cont. 2)

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 $1, \dots, K$.

observations	x_1	x_2	\dots	x_N
labels (unseen)	y_1	y_2	\dots	y_N

Then we would be able to use ML estimation recursively.

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 (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 Expectation-Maximisation Algorithm

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.

Limitation of ML-trained GMM

- Local optimum problem (the EM algorithm does not guarantee the global optimum)

- No mechanism for determining K (# of mixtures).

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

$$\gamma_k = \sum_{n=1}^N \gamma_k^n, \quad \gamma_k^n = \frac{w_k^{(t)} p(x_n|\lambda_k^{(t)})}{p(x_n|\Lambda^{(t)})}$$

$$Q_\pi(\Lambda^{(t)}, \lambda_k^{(t+1)}) = \sum_{n=1}^N \gamma_k^n \log p(x_n|\lambda_k^{(t+1)})$$

$$w_k^{(t+1)} = \gamma_k / \sum_{k=1}^K \gamma_k = \gamma_k / N$$

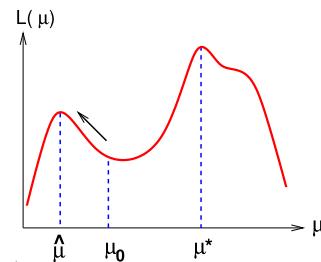
$$m_k^{(t+1)} = \sum_{n=1}^N \gamma_k^n x_n / \sum_{n=1}^N \gamma_k^n, \quad \sigma_k^{(t+1)} = \sum_{n=1}^N \gamma_k^n (x - m_k^{(t)})^2 / \sum_{n=1}^N \gamma_k^n$$

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 (cont. 2)

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 $k \leftarrow k * 2$.)

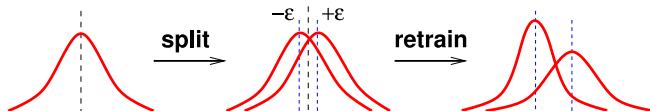
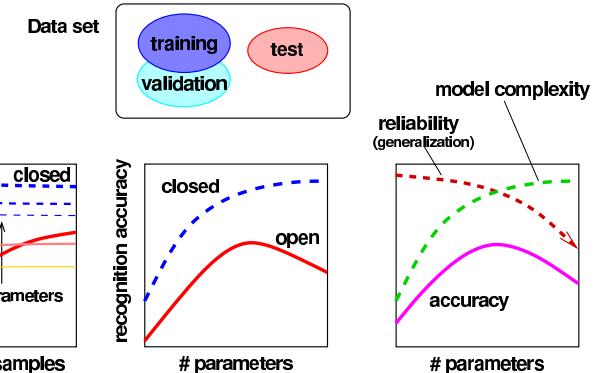


Fig. Mixture splitting

Remaining problems with GMM (cont. 2)



Remaining problems with GMM

- 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

Maximum a posterior (MAP) estimation

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

$$\Lambda : \text{a set of parameters to estimate}$$

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 Λ
 $p(\Lambda|X)$: 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

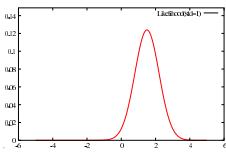
$$\hat{p}(\mathbf{x}|X) = p(\mathbf{x}|X, \Lambda^*), \quad \Lambda^* = \begin{cases} \max_{\Lambda} p(X|\Lambda), & \text{MLE} \\ \max_{\Lambda} p(X|\Lambda)p(\Lambda), & \text{MAP} \end{cases}$$

assuming parameter Λ is fixed but unknown.

- interval estimation (Bayesian estimation)

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

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



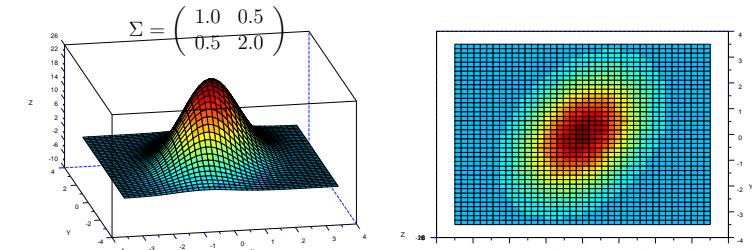
assuming the parameters Λ are random variables with a prior distribution $p(\Lambda)$

Multivariate Gaussian pdf_(cont. 2)

mean vector: $\mu = (\mu_1, \mu_2)^t$
covariance matrix: $\Sigma = \begin{pmatrix} \sigma_{11} & \sigma_{12} \\ \sigma_{21} & \sigma_{22} \end{pmatrix}$

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

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



Multivariate Gaussian pdf

One variable

data: $\mathbf{x}_n = (x_n)$
 $\mu = \frac{1}{N} \sum_{i=1}^N x_i, \quad \sigma^2 = \frac{1}{N} \sum_{i=1}^N (x_i - \mu)^2$

Two variables

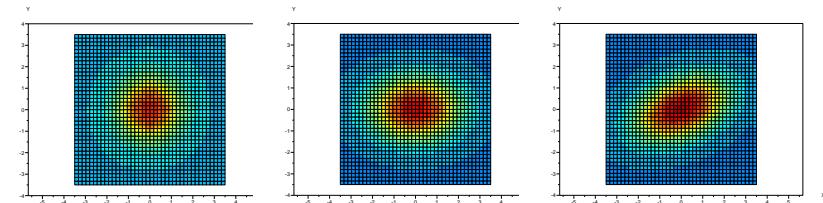
$$\begin{aligned} \mathbf{x}_n &= (x_{1n}, x_{2n})^t = \begin{pmatrix} x_{1n} \\ x_{2n} \end{pmatrix}, \quad \mu = \begin{pmatrix} \mu_1 \\ \mu_2 \end{pmatrix}, \quad \Sigma = \begin{pmatrix} \sigma_{11} & \sigma_{12} \\ \sigma_{21} & \sigma_{22} \end{pmatrix} \\ \mu_1 &= \frac{1}{N} \sum_{i=1}^N x_{1i},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

Multivariate Gaussian pdf_(cont. 3)

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

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

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

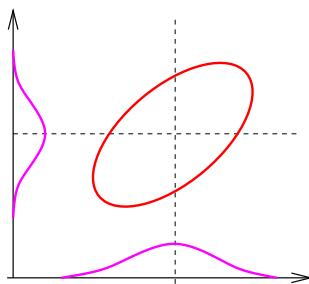


x_1 and x_2 are called **uncorrelated** if $\sigma_{12} = 0$.

Gaussian pdf with diagonal covariances

Gaussian pdf
with a full-covariance

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

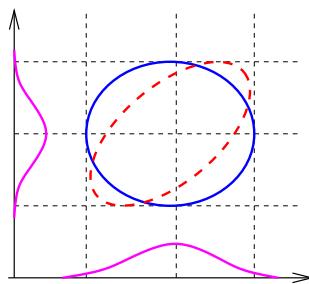


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Gaussian pdf
with diagonal-covariances

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

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



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Gaussian pdf with diagonal covariances (cont. 3)

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

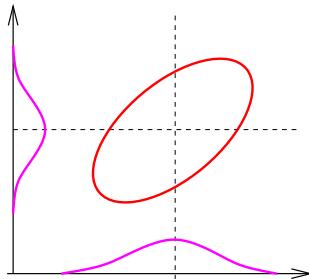
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Gaussian pdf with diagonal covariances (cont. 2)

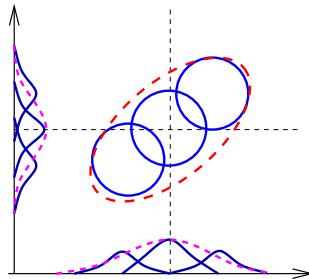
GMM with diag. cova. can model dependencies (correlations) !

Gaussian pdf
with full-covariance



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GMM pdf
with diagonal-covariance



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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

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

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.

k-fold CV	subsets
2-fold CV	Set2-0, Set2-1
5-fold CV	Set5-0, Set5-1, Set5-2, Set5-3, Set5-4

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

■ 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

Maximum Likelihood Estimation (MLE)

- http://en.wikipedia.org/wiki/Maximum_likelihood

GMM

- http://en.wikipedia.org/wiki/Mixture_model

EM algorithm

- http://en.wikipedia.org/wiki/Expectation-maximization_algorithm