## 02457 Non-linear signal processing

#### 2017 - Lecture 11



## **Technical University of Denmark**DTU Compute, Kgs Lyngby, Denmark



#### **Outline lecture 11**

#### Kernel methods

Dual representation, kernel matrices

$$(\mathbf{x}_m \approx \mathbf{x}_n) \Rightarrow (\mathbf{t}_m \approx \mathbf{t}_n)$$

#### Gaussian process prior

Smoothness ~ output correlation between input neighbors

#### **Support Vector Machines**

Sparse classifier

Sequence data — Ex Audio signals, Music and Speech, Speech production, Symbolic representations

#### Markov models

- Definition, properties, stationary distribution
- Sequencial Likelihood function and ML, MAP estimation

#### Dynamic estimators for non-stationary data

- Stochastic gradient
- Evaluation with non-stationarity



## Kernel representation of linear model d > N

Assume we have a high-dimensional data set with input variables **x** in d-dimensional space and d > N

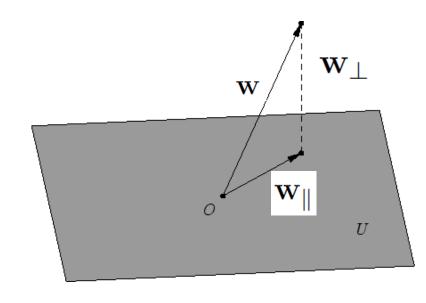
$$D = \{(t_1, \mathbf{x}_1), (t_2, \mathbf{x}_2), ..., (t_N, \mathbf{x}_N)\}\$$

Linear model fitted from least squares

$$E(\mathbf{w}) = \sum_{n=1}^{N} (t_n - \sum_{j=1}^{d+1} w_j x_{j,n})^2$$

$$\mathbf{w} = \mathbf{w}_{\perp} + \mathbf{w}_{\parallel}$$

$$\mathbf{w}_{\parallel} = \sum_{n=1}^{N} a_n \mathbf{x}_n, \qquad \mathbf{w}_{\perp}^{\top} \mathbf{x}_n = 0$$



## Kernel representation in the linear model

High-dimensional data d > N 
$$D = \{(t_1, \mathbf{x}_1), (t_2, \mathbf{x}_2), ..., (t_N, \mathbf{x}_N)\}$$
 
$$\mathbf{w} = \mathbf{w}_\perp + \mathbf{w}_\parallel \qquad \mathbf{w}_\parallel = \sum_{n=1}^N a_n \mathbf{x}_n, \quad \mathbf{w}_\perp^\top \mathbf{x}_n = 0$$

Linear model fitted from least squares

$$E(\mathbf{w}) = \sum_{n=1}^{N} (t_n - \mathbf{w}^{\top} \mathbf{x}_n)^2 = \sum_{n=1}^{N} (t_n - \mathbf{w}_{\parallel}^{\top} \mathbf{x}_n)^2 = \sum_{n=1}^{N} (t_n - \sum_{m=1}^{N} a_m \mathbf{x}_m^{\top} \mathbf{x}_n)^2$$
$$= \sum_{n=1}^{N} (t_n - \sum_{m=1}^{N} a_m K_{m,n})^2 = \sum_{n=1}^{N} (t_n - (\mathbf{a}^{\top} \mathbf{K})_n)^2.$$

$$(\mathbf{K})_{m,n} = K_{m,n} = \mathbf{x}_m^{\mathsf{T}} \mathbf{x}_n$$



#### Kernel methods

$$(\mathbf{K})_{m,n} = K_{m,n} = \mathbf{x}_m^{\mathsf{T}} \mathbf{x}_n$$

Note, we can ignore the component of the weight vector which is orthogonal to the subspace spanned by the data

Costfunction (likelihood function) is "blind" to this subspace

We can reduce the fitting problem to estimation of an N-dimensional vector (a) forget about high-dim data vectors x and just keep the kernel matrix K



## Kernel methods: Implict features

Assume our model is based on a feature representation –

$$\mathbf{x} \mapsto \phi(\mathbf{x})$$

with the kernel trick there is no limitations on the dimensionality

$$E(\mathbf{w}) = \sum_{n=1}^{N} (t_n - \mathbf{w}^{\mathsf{T}} \boldsymbol{\phi}(\mathbf{x}_n))^2 = \sum_{n=1}^{N} (t_n - (\mathbf{a}^{\mathsf{T}} \mathbf{K})_n)^2$$

$$(\mathbf{K})_{m,n} = K_{m,n} = \boldsymbol{\phi}(\mathbf{x}_m)^{\mathsf{T}} \boldsymbol{\phi}(\mathbf{x}_n)$$



## Kernel methods: Implict features

Now if we postulate a kernel function K(x,x') then we implicitly define a such feature representation

$$(\mathbf{K})_{m,n} = K_{m,n} = \boldsymbol{\phi}(\mathbf{x}_m)^{\top} \boldsymbol{\phi}(\mathbf{x}_n)$$

When can a symmetric matrix generated from a function  $k(\mathbf{x}_m, \mathbf{x}_n)$  be reconstructed as the inner products?

If any subset of points give rise to a positive definite matrix.



#### Techniques for Constructing New Kernels.

Given valid kernels  $k_1(\mathbf{x}, \mathbf{x}')$  and  $k_2(\mathbf{x}, \mathbf{x}')$ , the following new kernels will also be valid:

$$k(\mathbf{x}, \mathbf{x}') = ck_1(\mathbf{x}, \mathbf{x}') \tag{6.13}$$

$$k(\mathbf{x}, \mathbf{x}') = f(\mathbf{x})k_1(\mathbf{x}, \mathbf{x}')f(\mathbf{x}') \tag{6.14}$$

$$k(\mathbf{x}, \mathbf{x}') = q(k_1(\mathbf{x}, \mathbf{x}')) \tag{6.15}$$

$$k(\mathbf{x}, \mathbf{x}') = \exp(k_1(\mathbf{x}, \mathbf{x}')) \tag{6.16}$$

$$k(\mathbf{x}, \mathbf{x}') = k_1(\mathbf{x}, \mathbf{x}') + k_2(\mathbf{x}, \mathbf{x}')$$
 (6.17)

$$k(\mathbf{x}, \mathbf{x}') = k_1(\mathbf{x}, \mathbf{x}')k_2(\mathbf{x}, \mathbf{x}') \tag{6.18}$$

$$k(\mathbf{x}, \mathbf{x}') = k_3(\phi(\mathbf{x}), \phi(\mathbf{x}')) \tag{6.19}$$

$$k(\mathbf{x}, \mathbf{x}') = \mathbf{x}^{\mathrm{T}} \mathbf{A} \mathbf{x}' \tag{6.20}$$

$$k(\mathbf{x}, \mathbf{x}') = k_a(\mathbf{x}_a, \mathbf{x}'_a) + k_b(\mathbf{x}_b, \mathbf{x}'_b) \tag{6.21}$$

$$k(\mathbf{x}, \mathbf{x}') = k_a(\mathbf{x}_a, \mathbf{x}'_a)k_b(\mathbf{x}_b, \mathbf{x}'_b) \tag{6.22}$$

where c > 0 is a constant,  $f(\cdot)$  is any function,  $q(\cdot)$  is a polynomial with nonnegative coefficients,  $\phi(\mathbf{x})$  is a function from  $\mathbf{x}$  to  $\mathbb{R}^M$ ,  $k_3(\cdot, \cdot)$  is a valid kernel in  $\mathbb{R}^M$ ,  $\mathbf{A}$  is a symmetric positive semidefinite matrix,  $\mathbf{x}_a$  and  $\mathbf{x}_b$  are variables (not necessarily disjoint) with  $\mathbf{x} = (\mathbf{x}_a, \mathbf{x}_b)$ , and  $k_a$  and  $k_b$  are valid kernel functions over their respective spaces.



# Gaussian kernel, aka squared exponentil aka radial basis function kernel

$$k(\mathbf{x}, \mathbf{x}') = \exp\left(-\|\mathbf{x} - \mathbf{x}'\|^2 / 2\sigma^2\right)$$

#### Is it an ok kernel?

$$\|\mathbf{x} - \mathbf{x}'\|^2 = \mathbf{x}^T \mathbf{x} + (\mathbf{x}')^T \mathbf{x}' - 2\mathbf{x}^T \mathbf{x}'$$

$$k(\mathbf{x}, \mathbf{x}') = \exp(-\mathbf{x}^{\mathrm{T}}\mathbf{x}/2\sigma^{2}) \exp(\mathbf{x}^{\mathrm{T}}\mathbf{x}'/\sigma^{2}) \exp(-(\mathbf{x}')^{\mathrm{T}}\mathbf{x}'/2\sigma^{2})$$

$$k(\mathbf{x}, \mathbf{x}') = f(\mathbf{x})k_1(\mathbf{x}, \mathbf{x}')f(\mathbf{x}')$$
 (6.14)

$$k(\mathbf{x}, \mathbf{x}') = \exp(k_1(\mathbf{x}, \mathbf{x}')) \tag{6.16}$$



## Kernel methods for supervised learning

The general idea of kernel representations for supervised learning is we can implement smoothness as:

Similarity in input => similarity in output

$$(\mathbf{x}_m \approx \mathbf{x}_n) \Rightarrow (\mathbf{t}_m \approx \mathbf{t}_n)$$

## Gaussian processes for function approximation

$$(\mathbf{x}_m \approx \mathbf{x}_n) \Rightarrow (\mathbf{t}_m \approx \mathbf{t}_n)$$

In the Gaussian process model the similarity is implemented in a probabilistic setting

For additive noise model  $t(\mathbf{x}) = y(\mathbf{x}) + e$ , we can represent the similarity as an assumed Gaussian distribution of the function values for a general set of inputs

$$cov(y_1, ..., y_N) = \mathbf{K}, \quad p(\mathbf{y}|\mathbf{K}) = \frac{1}{|2\pi\mathbf{K}|^{\frac{1}{2}}} \exp(-\frac{1}{2}\mathbf{y}^{\mathsf{T}}\mathbf{K}^{-1}\mathbf{y})$$



## Kernel methods for supervised learning

The Gaussian distribution of the target then follows

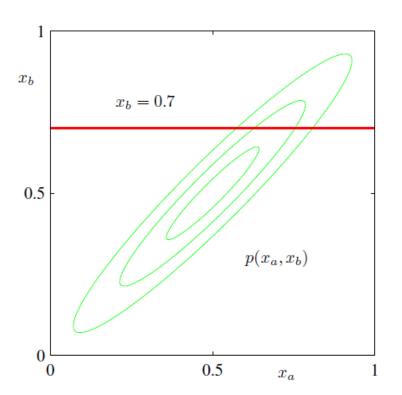
$$cov(t_1, ..., t_N) \equiv \mathbf{C} = \mathbf{K} + \beta^{-1}\mathbf{1}$$

Predictive distribution

$$p(\mathbf{t}_{\text{test}}|\mathbf{t}_{\text{train}}, \mathbf{C}_{\text{train},\text{test}}) = \frac{p(\mathbf{t}_{\text{test}}, \mathbf{t}_{\text{train}}|\mathbf{C}_{\text{train},\text{test}})}{p(\mathbf{t}_{\text{train}}|\mathbf{C}_{\text{train}})}$$



## Gaussian conditioning



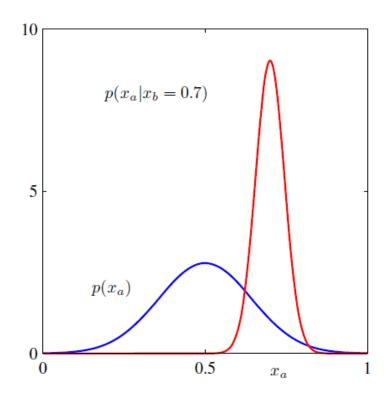


Figure 2.9 The plot on the left shows the contours of a Gaussian distribution  $p(x_a, x_b)$  over two variables, and the plot on the right shows the marginal distribution  $p(x_a)$  (blue curve) and the conditional distribution  $p(x_a|x_b)$  for  $x_b = 0.7$  (red curve).

## Kernel methods for supervised learning

## Rules for conditioning

(Bishop Eq. (281-82))

$$\mu_{\mathrm{test|train}} = \mathrm{C}_{\mathrm{test,train}} \mathrm{C}_{\mathrm{train}}^{-1} \mathrm{t}_{\mathrm{train}}$$

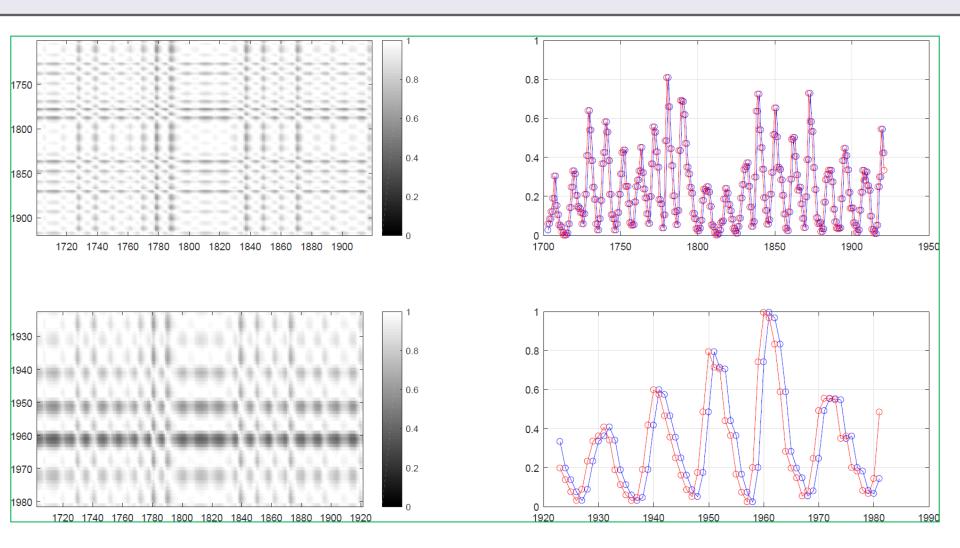
$$\mathbf{C}_{test|train} \ = \ \mathbf{C}_{test} - \mathbf{C}_{test,train} \mathbf{C}_{train}^{-1} \mathbf{C}_{train,test}$$

$$\hat{t}_m = (\boldsymbol{\mu}_{\text{test|train}})_m$$

$$\operatorname{std}(t_m) = \sqrt{(\mathbf{C}_{\operatorname{test|train}})_{m,m}} = \sqrt{(\mathbf{C}_{\operatorname{test}})_{m,m} - (\mathbf{C}_{\operatorname{test,train}}\mathbf{C}_{\operatorname{train}}^{-1}\mathbf{C}_{\operatorname{train,test}})_{m,m}}$$

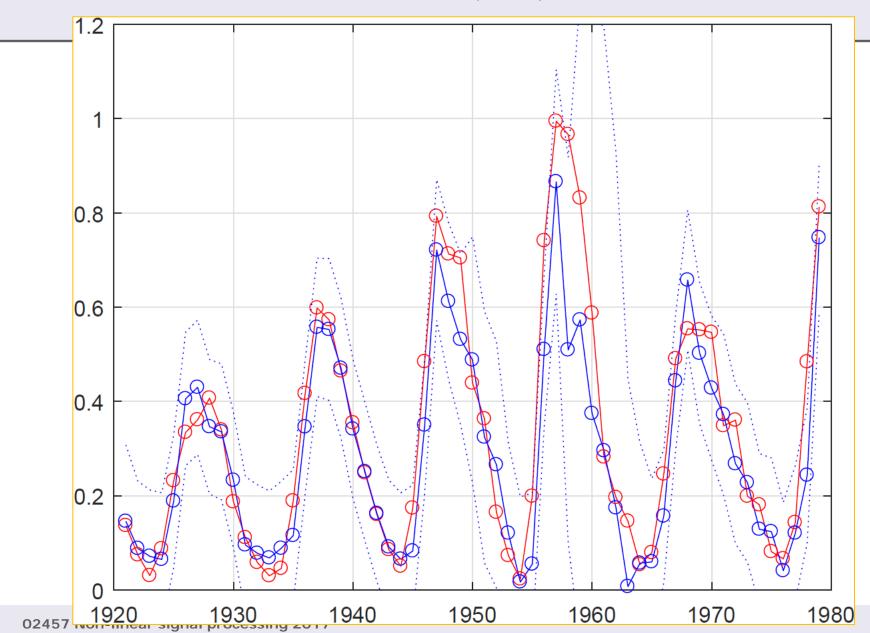


## Sun spot Ctrain, Ctest, train





## Confidence interval $(2*\sigma)$



## A simple alternative derivation of the GP prediction (Bishop sec. 6.1)

$$J(\mathbf{w}) = \frac{1}{2} \sum_{n=1}^{N} \left\{ \mathbf{w}^{\mathrm{T}} \phi(\mathbf{x}_n) - t_n \right\}^2 + \frac{\lambda}{2} \mathbf{w}^{\mathrm{T}} \mathbf{w}$$

$$\mathbf{w} = -\frac{1}{\lambda} \sum_{n=1}^{N} \left\{ \mathbf{w}^{\mathrm{T}} \phi(\mathbf{x}_n) - t_n \right\} \phi(\mathbf{x}_n) = \sum_{n=1}^{N} a_n \phi(\mathbf{x}_n) = \mathbf{\Phi}^{\mathrm{T}} \mathbf{a}$$

$$J(\mathbf{a}) = \frac{1}{2}\mathbf{a}^{\mathrm{T}}\boldsymbol{\Phi}\boldsymbol{\Phi}^{\mathrm{T}}\boldsymbol{\Phi}\boldsymbol{\Phi}^{\mathrm{T}}\mathbf{a} - \mathbf{a}^{\mathrm{T}}\boldsymbol{\Phi}\boldsymbol{\Phi}^{\mathrm{T}}\mathbf{t} + \frac{1}{2}\mathbf{t}^{\mathrm{T}}\mathbf{t} + \frac{\lambda}{2}\mathbf{a}^{\mathrm{T}}\boldsymbol{\Phi}\boldsymbol{\Phi}^{\mathrm{T}}\mathbf{a}$$

$$J(\mathbf{a}) = \frac{1}{2} \mathbf{a}^{\mathrm{T}} \mathbf{K} \mathbf{K} \mathbf{a} - \mathbf{a}^{\mathrm{T}} \mathbf{K} \mathbf{t} + \frac{1}{2} \mathbf{t}^{\mathrm{T}} \mathbf{t} + \frac{\lambda}{2} \mathbf{a}^{\mathrm{T}} \mathbf{K} \mathbf{a}. \qquad => \qquad \mathbf{a} = (\mathbf{K} + \lambda \mathbf{I}_N)^{-1} \mathbf{t}.$$

$$y(\mathbf{x}) = \mathbf{w}^{\mathrm{T}} \phi(\mathbf{x}) = \mathbf{a}^{\mathrm{T}} \mathbf{\Phi} \phi(\mathbf{x}) = \mathbf{k}(\mathbf{x})^{\mathrm{T}} (\mathbf{K} + \lambda \mathbf{I}_{N})^{-1} \mathbf{t}$$



## Support vector machines

The general idea of kernel representations for supervised learning is that similarity in input should lead to similarity in output

For support vectors this is assumed for labels

$$(\mathbf{x}_m \approx \mathbf{x}_n) \Rightarrow (\mathbf{t}_m \approx \mathbf{t}_n)$$

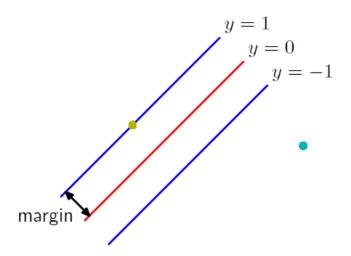
Use kernel to implement similarity and linear discriminant to make decision

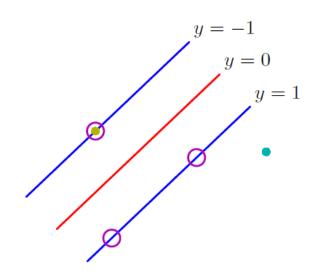


## Maximum margin principle

$$y(\mathbf{x}) = \mathbf{w}^{\mathrm{T}} \phi(\mathbf{x}) + b$$

Margin = 
$$|y(\mathbf{x})|/||\mathbf{w}||$$

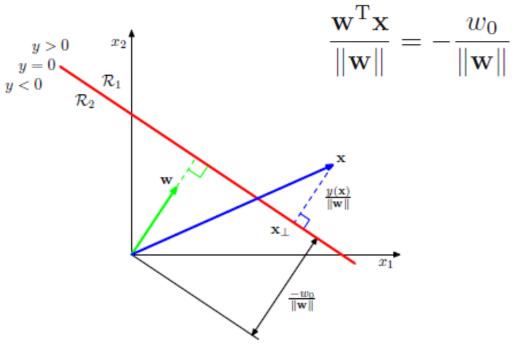




**Figure 7.1** The margin is defined as the perpendicular distance between the decision boundary and the closest of the data points, as shown on the left figure. Maximizing the margin leads to a particular choice of decision boundary, as shown on the right. The location of this boundary is determined by a subset of the data points, known as support vectors, which are indicated by the circles.

## Geometry of linear discriminant

Figure 4.1 Illustration of the geometry of a linear discriminant function in two dimensions. The decision surface, shown in red, is perpendicular to  $\mathbf{w},$  and its displacement from the origin is controlled by the bias parameter  $w_0.$  Also, the signed orthogonal distance of a general point  $\mathbf{x}$  from the decision surface is given by  $y(\mathbf{x})/\|\mathbf{w}\|.$ 



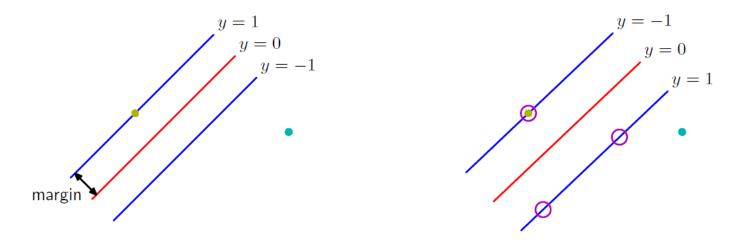
$$\mathbf{x} = \mathbf{x}_{\perp} + r \frac{\mathbf{w}}{\|\mathbf{w}\|}.$$

$$y(\mathbf{x}) = \mathbf{w}^{\mathrm{T}} \phi(\mathbf{x}) + b$$

## Maximum margin principle

Margin = 
$$|y(\mathbf{x})|/||\mathbf{w}||$$
 
$$\frac{t_n y(\mathbf{x}_n)}{||\mathbf{w}||} = \frac{t_n(\mathbf{w}^T \phi(\mathbf{x}_n) + b)}{||\mathbf{w}||}$$

$$\underset{\mathbf{w},b}{\operatorname{arg\,max}} \left\{ \frac{1}{\|\mathbf{w}\|} \min_{n} \left[ t_n \left( \mathbf{w}^{\mathrm{T}} \boldsymbol{\phi}(\mathbf{x}_n) + b \right) \right] \right\}$$



**Figure 7.1** The margin is defined as the perpendicular distance between the decision boundary and the closest of the data points, as shown on the left figure. Maximizing the margin leads to a particular choice of decision boundary, as shown on the right. The location of this boundary is determined by a subset of the data points, known as support vectors, which are indicated by the circles.

#### Convex criteria

rescaling  $\mathbf{w} \to \kappa \mathbf{w}$  and  $b \to \kappa b$ 

$$t_n\left(\mathbf{w}^{\mathrm{T}}\phi(\mathbf{x}_n) + b\right) = 1$$
 For points closest to y=0

$$t_n\left(\mathbf{w}^{\mathrm{T}}\boldsymbol{\phi}(\mathbf{x}_n) + b\right) \geqslant 1$$
 For all points

$$\underset{\mathbf{w},b}{\operatorname{arg\,max}} \left\{ \frac{1}{\|\mathbf{w}\|} \min_{n} \left[ t_n \left( \mathbf{w}^{\mathrm{T}} \boldsymbol{\phi}(\mathbf{x}_n) + b \right) \right] \right\}$$

$$\underset{\mathbf{w},b}{\arg\min} \frac{1}{2} \|\mathbf{w}\|^2$$
 with the above constraint

#### **Convex optimization**

$$L(\mathbf{w}, b, \mathbf{a}) = \frac{1}{2} ||\mathbf{w}||^2 - \sum_{n=1}^{N} a_n \left\{ t_n(\mathbf{w}^{\mathrm{T}} \boldsymbol{\phi}(\mathbf{x}_n) + b) - 1 \right\}$$

$$\mathbf{w} = \sum_{n=1}^{N} a_n t_n \phi(\mathbf{x}_n)$$
$$0 = \sum_{n=1}^{N} a_n t_n.$$

#### Convex optimization = quadratic optimization w. constraint

Eliminating w and b from  $L(\mathbf{w}, b, \mathbf{a})$  using these conditions then gives the dual representation of the maximum margin problem in which we maximize

$$\widetilde{L}(\mathbf{a}) = \sum_{n=1}^{N} a_n - \frac{1}{2} \sum_{n=1}^{N} \sum_{m=1}^{N} a_n a_m t_n t_m k(\mathbf{x}_n, \mathbf{x}_m)$$
 (7.10)

with respect to a subject to the constraints

$$a_n \geqslant 0, \qquad n = 1, \dots, N, \tag{7.11}$$

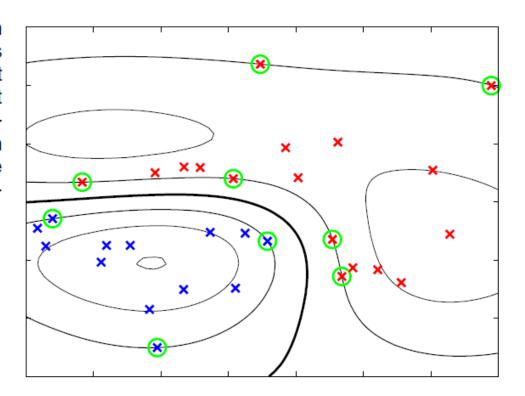
$$\sum_{n=1}^{N} a_n t_n = 0. (7.12)$$

$$k(\mathbf{x}, \mathbf{x}') = \phi(\mathbf{x})^{\mathrm{T}} \phi(\mathbf{x}')$$

$$\hat{t}_m = \operatorname{sign}\left(\sum_{n=1}^N a_n t_n K(\mathbf{x}_m, \mathbf{x}_n) + b\right)$$



Figure 7.2 Example of synthetic data from two classes in two dimensions showing contours of constant  $y(\mathbf{x})$  obtained from a support vector machine having a Gaussian kernel function. Also shown are the decision boundary, the margin boundaries, and the support vectors.



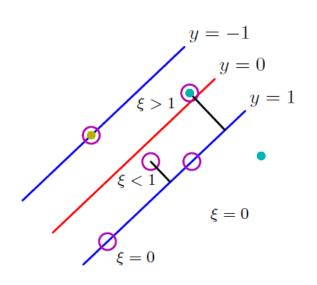
## Slack variables for noisy data (non-separable)

**Figure 7.3** Illustration of the slack variables  $\xi_n \geqslant 0$ . Data points with circles around them are support vectors.

$$t_n y(\mathbf{x}_n) \geqslant 1 - \xi_n,$$

Cost associated with slack

$$C\sum_{n=1}^{N} \xi_n + \frac{1}{2} \|\mathbf{w}\|^2$$



Augmented Lagrange function (a, µ parameters positive)

$$L(\mathbf{w}, b, \mathbf{a}) = \frac{1}{2} \|\mathbf{w}\|^2 + C \sum_{n=1}^{N} \xi_n - \sum_{n=1}^{N} a_n \{t_n y(\mathbf{x}_n) - 1 + \xi_n\} - \sum_{n=1}^{N} \mu_n \xi_n$$

$$\frac{\partial L}{\partial \xi_n} = 0 \quad \Rightarrow \quad a_n = C - \mu_n \qquad = > \qquad 0 \leqslant a_n \leqslant C$$

$$0 \leqslant a_n \leqslant C$$

#### Slack variables for noisy data (non-separable)

#### Modified optimization problem

$$\tilde{L}(\mathbf{a}) = \sum_{n=1}^{N} a_n - \sum_{n,m=1}^{N} a_n a_m t_n t_m K(\mathbf{x}_n, \mathbf{x}_m),$$

$$0 \le a_n \le C,$$

$$\sum_{n=1}^{N} a_n t_n = 0,$$

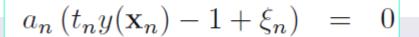
$$b = \frac{1}{|M|} \sum_{n \in M} (t_n - \sum_{m \in S} a_m t_m K(\mathbf{x}_n, \mathbf{x}_m)).$$

Interpretation of coefficients:

a > 0: support vectors, enters classification

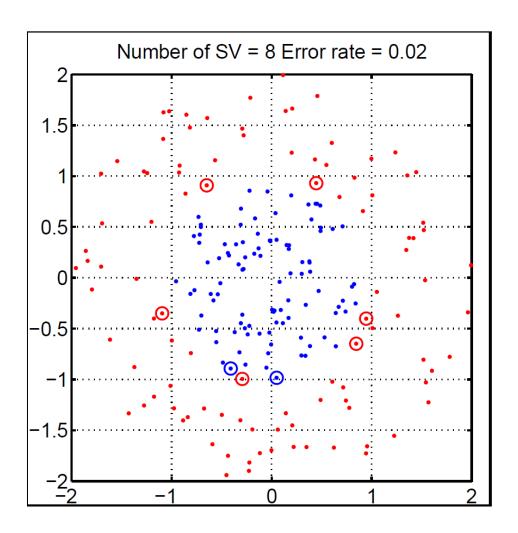
0< a < C: on the boundary

a = C: slack "activated" => inside margin; misclass if  $\xi > 1$ 





## Simple 2D synthetic data



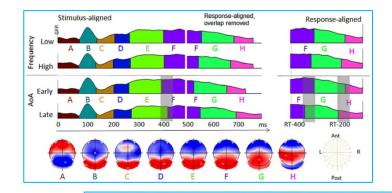


#### Dynamical systems represented as symbolic sequences

Many real life applications are based on symbolic, label

sequence based representations

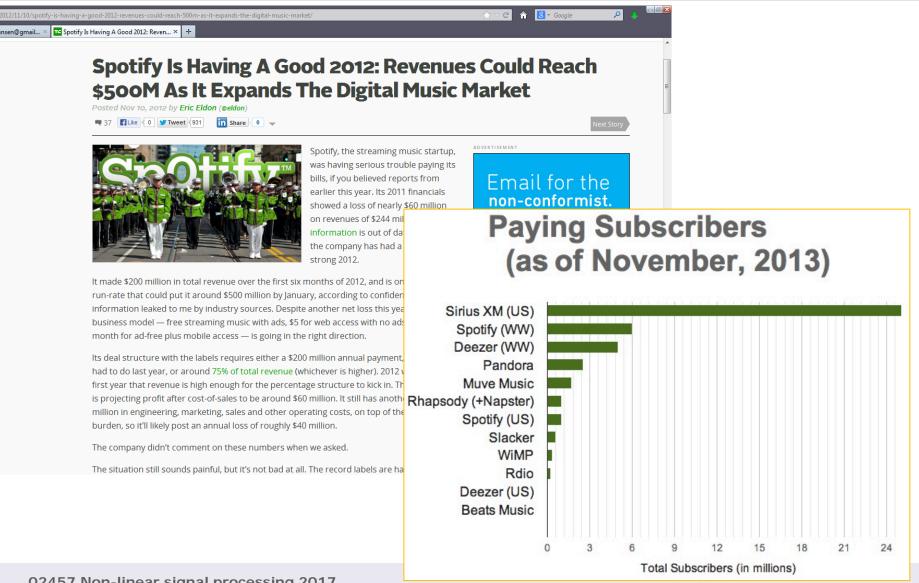
Speech, phonemes
Gesture, activity recognition
Text, words, topics
Mobility, stop locations
EEG, micro-states, meso-states)







## **Spotify**







# WinAmp demo project





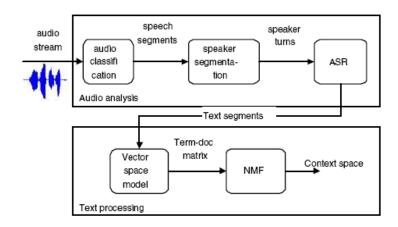


Lehn-Schiøler, Garcias-Arena, Petersen, Hansen: ISMIR 2006 (Oct 9, 2006)

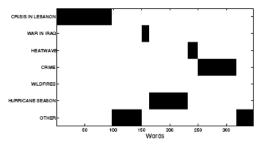
#### CASTSEARCH - CONTEXT BASED SPEECH DOCUMENT RETRIEVAL

Lasse Lohilahti Mølgaard, Kasper Winther Jørgensen, and Lars Kai Hansen

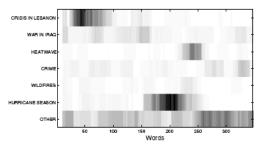
Informatics and Mathematical Modelling
Technical University of Denmark Richard Petersens Plads
Building 321, DK-2800 Kongens Lyngby, Denmark



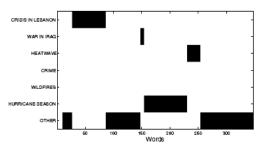
**Fig. 1.** The system setup. The audio stream is first processed using audio segmentation. Segments are then using an automatic speech recognition (ASR) system to produce text segments. The text is then processed using a vector representation of text and apply nonnegative matrix factorization (NMF) to find a topic space.



(a) Manual segmentation.



(b)  $p(k|d^*)$  for each context. Black means high probability.



(c) The segmentation based on p(k|d\*).

Fig. 3. Figure 3(a) shows the manual segmentation of the news show into 7 classes. Figure 3(b) shows the distribution  $p(k|d^*)$  used to do the actual segmentation shown in figure 3(c). The NMF-segmentation is in general consistent with the manual segmentation. Though, the segment that is manually segmented as 'crime' is labeled 'other' by the NMF-segmentation



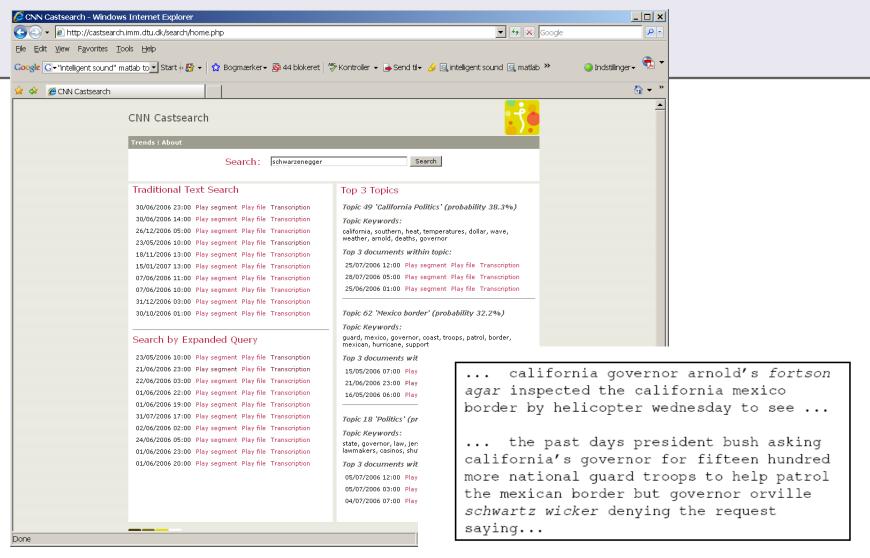
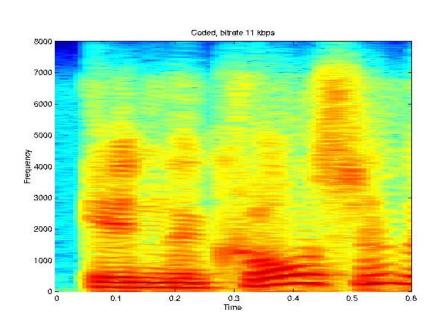
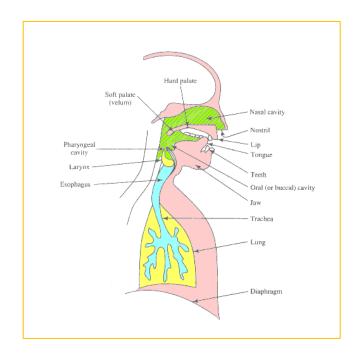


Fig. 2. Two examples of the retrieved text for a query on 'schwarzenegger'.

## castsearch.imm.dtu.dk

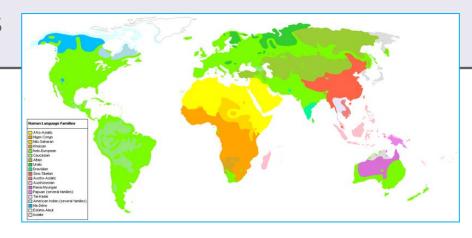
## Sequence model example: Audio DSP







#### **SPEECH** –some definitions



- Speech signals are sequences of sounds
- The basic sounds and the transitions between them serve as symbolic representation of information - semantics
- The arrangement of these sounds (symbols) is governed by the rules of *language*
- The study of these rules and their implications in human communication is called *linguistics*
- The study of and classification of the sounds of speech is called *phonetics*



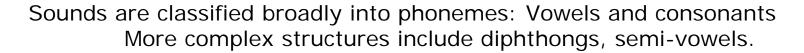
#### SPEECH PRODUCTION

Speech is produced by the human vocal tract

The vocal tract is excited either by short burst of periodic stimulus or white noise

Voiced sounds are produced by an airflow through tight vocal cords.

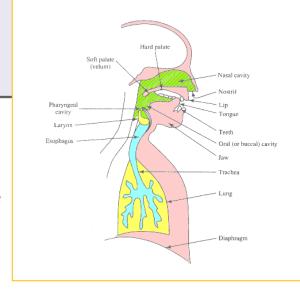
Unvoiced sounds are produced by turbulent flow.

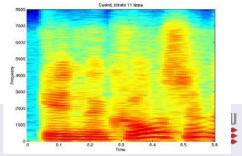


Formants are peaks in the power spectrum modulation (envelopes).

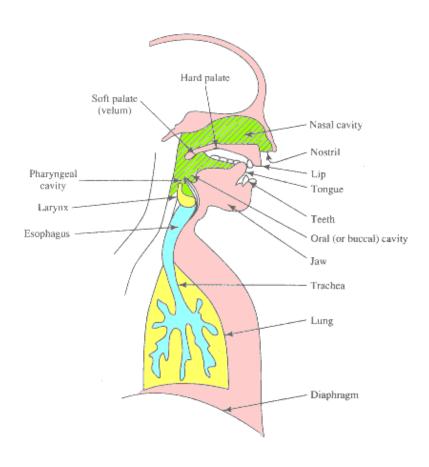
Frequencies in the range:

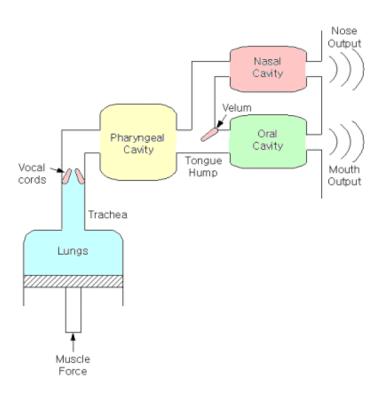
F1 (270-730 Hz), F2 (840-2290 Hz), and F3 (1690-3010 Hz).





# Modeling speech

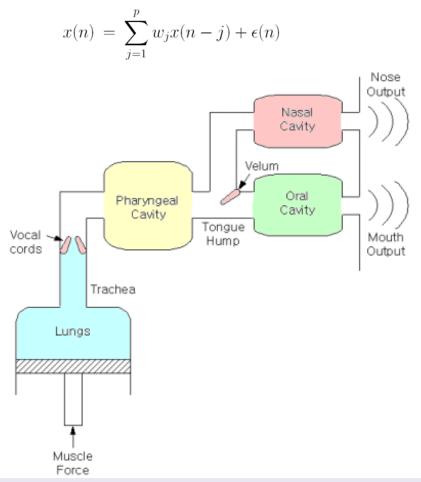




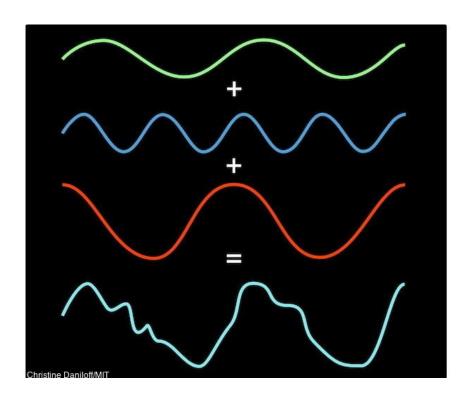


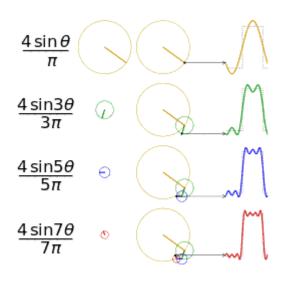
# The speech signal is modelled as a linear time variant system excited by a high-frequency signal (periodic og random)

• Linear system model (IIR model)



#### Fourier transform



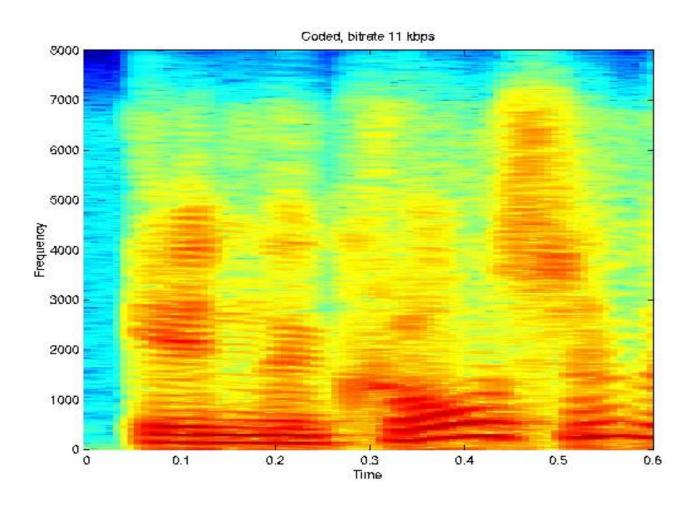


https://en.wikipedia.org/wiki/Fourier\_series

http://en.wikipedia.org/wiki/Discrete\_Fourier\_transform



# Speech spectrogram





#### CEPSTRAL COEFFICIENTS

Linear system model (IIR model)

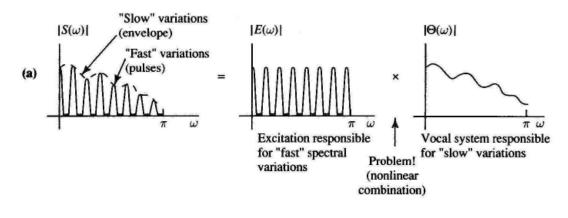
$$s(n) = \sum_{j=1}^{p} w(j) \cdot s(n-j) + \epsilon(n)$$
  
$$s(n) = \sum_{i}^{p} \theta(i) \cdot \epsilon(n-i)$$

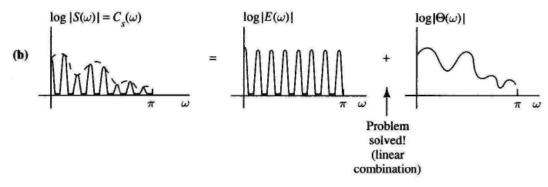
In the Fourier domain

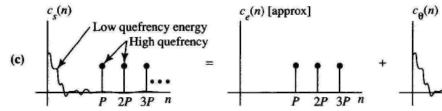
$$S(\omega) = \Theta(\omega) \cdot E(\omega)$$
$$\log |S(\omega)| = \log |\Theta(\omega)| + \log |E(\omega)|$$

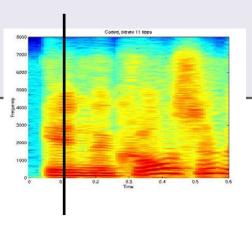


# **Cepstral liftering**





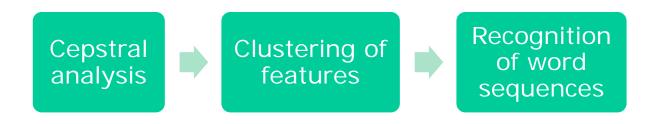




## Phoneme grouping and sequence recognition

Provides discrete symbols for identification of word.

The audio stream is segmented into a string of symbols (~phonemes), by assigning windows to most likely phoneme using K-means clustering.







A string of symbols 'y<sub>n</sub>' is a Markov chain if

$$p(y_n \mid y_1, y_2, ..., y_{n-2}, y_{n-1}) = p(y_n \mid y_{n-1}, ..., y_{n-2}, y_{n-1})$$

A string of symbols is 1'st order Markovian if

$$p(y_n | y_1, y_2, ..., y_{n-2}, y_{n-1}) = p(y_n | y_{n-1})$$



#### SIMPLE MARKOV MODEL

Let  $y_n$  be a sequence of symbols with K states

Let  $a_{j,j'}$  be the probability of jumping from j to j'

i.e., 
$$a_{j,j'} = p(y_n = j' | y_{n-1} = j)$$

Matrix  $a_{j,j'}$  is a stochastic matrix  $\Sigma_{j'}$   $a_{j,j'} = 1$ 

a can be estimated by maximum likelihood



#### Markov chain estimation by maximum likelihood

The likelihood can be written, introducing the observed number of transitions  $j = j' : n_{i,i'}$ 

$$P(\{y_n\}|a) = P(y_1) \prod_{n=2}^{N} P(y_n|y_{n-1}, a)$$
$$= P(y_1) \prod_{j,j'} (a_{j,j'})^{n_{j,j'}}$$

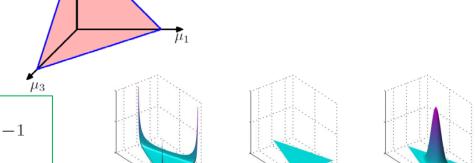
**Estimator** 

$$\widehat{a}_{j,j'} = \frac{n_{j,j'}}{\sum_{j''} n_{j,j''}}$$



### MAP estimate with Dirichlet priors on aj,:

The Dirichlet distribution over three variables  $\mu_1, \mu_2, \mu_3$  is confined to a simplex (a bounded linear manifold) of the form shown, as a consequence of the constraints  $0 \leqslant \mu_k \leqslant 1$  and  $\sum_k \mu_k = 1$ .



$$\operatorname{Dir}(\boldsymbol{\mu}|\boldsymbol{\alpha}) = \frac{\Gamma(\alpha_0)}{\Gamma(\alpha_1)\cdots\Gamma(\alpha_K)} \prod_{k=1}^K \mu_k^{\alpha_k - 1}$$

Figure 2.5 Plots of the Dirichlet distribution over three variables, where the two horizontal axes are coordinates in the plane of the simplex and the vertical axis corresponds to the value of the density. Here  $\{\alpha_k\}=0.1$  on the left plot,  $\{\alpha_k\}=1$  in the centre plot, and  $\{\alpha_k\}=10$  in the right plot.

$$\widehat{a}_{j,j'} = \frac{n_{j,j'} + \alpha_{j,j'} - 1}{\sum_{j''} n_{j,j''} + \alpha_{j,j''} - 1}$$
$$= \frac{n_{j,j'} + \alpha_{j,j'} - 1}{\sum_{j''} n_{j,j''} + \alpha_{j,j''} - 1}$$

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#### Analysis of Markov chains - The ensemble picture

Consider a large number of 'parallel' Markov chains  $(y^i)_n$ Each chain makes random moves according to a common a-matrix (stationarity) The probability over states, at time step n, obeys the Kolmogorov Chapman equation

$$P_{n+1}(j') = \sum_{j=1}^{K} P_n(j) a_{j,j'}$$
  $\mathbf{P}_{n+1} = \mathbf{P}_n \mathbf{A}$ 

The stationary distribution (Perron-Frobenius theorem, criterion for uniqueness) is a (left) eigenvector, with unit eigenvalue

$$P_*(j') = \sum_{j=1}^K P_*(j)a_{j,j'}. \qquad \mathbf{P_*} = \mathbf{P_*A}$$



#### **Detection based on sequences**

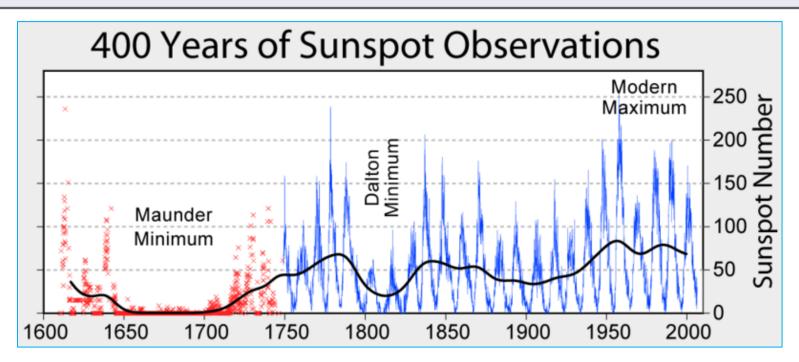
$$P(B|\{y_n\}) = \frac{P(\{y_n\}|B)P(B)}{P(\{y_n\})}.$$

$$P(\mathbf{A}_{B}|\{y_{n}\}) = \frac{P(\{y_{n}\}|\mathbf{A}_{B})P(\mathbf{A}_{B})}{P(\{y_{n}\})}.$$

$$P(\{y_n\}|a) = P(y_1) \prod_{n=2}^{N} P(y_n|y_{n-1}, a)$$
$$= P(y_1) \prod_{j,j'} (a_{j,j'})^{n_{j,j'}}$$



# What if data is non-stationary?



## Issues

How to track the changes in the model? Window based - assuming local stationarity Tracking slow changes

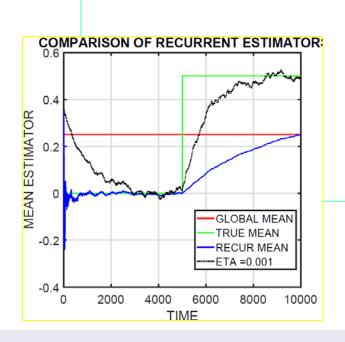
How to evaluate the error?

"Next sample error" - window



## Dynamic estimator of the mean

Dynamic updates for stream of data  $\{x_1, x_2, ..., x_N\}, \quad \mu = \frac{1}{N} \sum_{n=1}^{N} x_n$ 



$$\mu_{N} = \frac{1}{N} x_{N} + \frac{1}{N} \sum_{n=1}^{N-1} x_{n}$$

$$= \frac{1}{N} x_{N} + \frac{N-1}{N} \frac{1}{N-1} \sum_{n=1}^{N-1} x_{n}$$

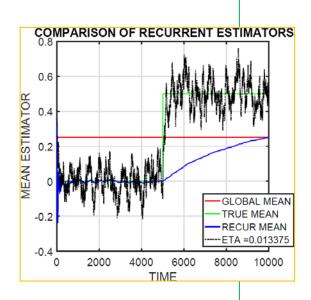
$$= \frac{1}{N} x_{N} + \frac{N-1}{N} \mu_{N-1}$$

$$= \mu_{N-1} + \frac{1}{N} (x_{N} - \mu_{N-1})$$

Non-stationarity, stochastic gradient

$$\mu_{N} = \mu_{N-1} - \eta \frac{\partial}{\partial \mu} \left[ \frac{1}{2} (x_{N} - \mu)^{2} \right]_{N-1}$$
$$= \mu_{N-1} + \eta (x_{N} - \mu_{N-1})$$

Difference equation has explicit solution



$$\mu_{N} = \sum_{n=1}^{N} \eta (1 - \eta)^{N-n} x_{n}$$

$$= \eta \sum_{n=1}^{N} \exp((N - n) \log(1 - \eta)) x_{n}$$

$$\approx \eta \sum_{n=1}^{N} \exp(-(N - n)\eta) x_{n}$$

$$\approx \eta \sum_{q=1}^{N} \exp(-q\eta) x_{N-q}$$

$$\approx \frac{1}{W} \sum_{q=1}^{W} x_{N-q}$$

#### Compare dynamic estimators in non-stationary data

