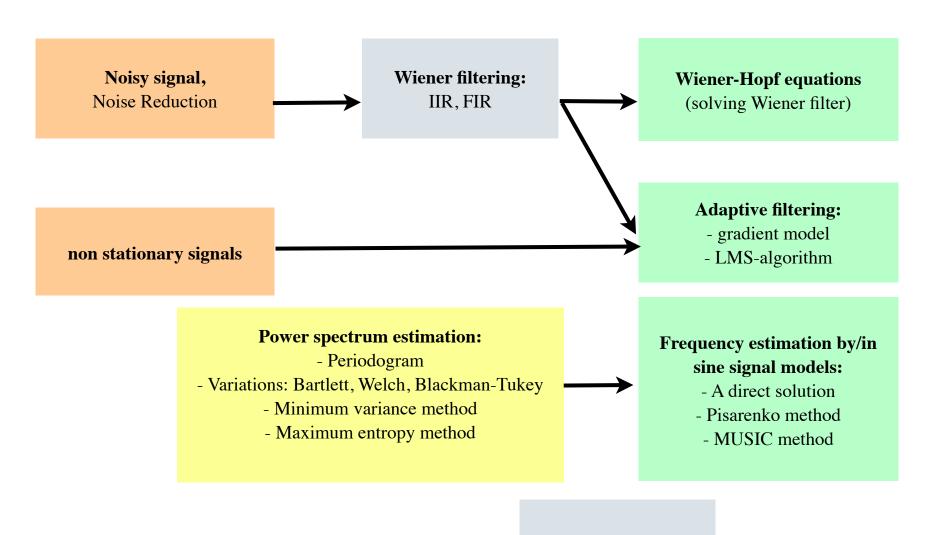


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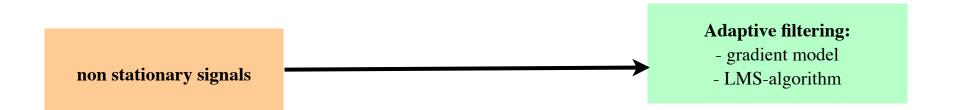
Adaptive filtering

#### Diagram of content of the final part of the course



**Estimation of variance:** 

#### Diagram of content of the final part of the course



#### **Estimation of variance:**

- Methods discussed so far require that the process is stationary (WSS) and ergodic
- Estimation can then be performed by averaging over time and using all observations  $x(0), x(1), \dots, x(N-1)$
- In many applications, stationarity can not always be assumed



- Examples:
- Speech signal: for example, the frequency and/or intensity of a vowel can change during the pronunciation
- For a signal measured from an electric motor, the frequency changes if the engine speed changes



- The default "process is non-stationary" is by itself far too weak
- We need more accurate assumptions, which are usually application-specific
- "Almost WSS": statistical properties change slowly over time
- Simple solution: look at the process in pieces and assume each piece is a WSS process



- Even if piecewise estimation were successful, estimated values change abruptly
- In many applications this may be an unwanted property
- However, if the parameters of interest do not change over time, then cutting the process may be sufficient



- Example: x(n) = -a(1,n) x(n-1) a(2) x(n-2) + v(n)
- a(1,n) is a function of time, but a(2) is constant
- Previous methods do not work, because a(1) does not remain constant
- Piecewise modeling provides several estimates of a(2), for which the average can be calculated
- demo: nonwss.R



- In general, we are interested in all the parameters
- In addition, often you want to track slow changes smoothly so cutting is not a satisfactory solution
- Consider, therefore, how to deal with processes which are assumed to be non stationary, but changes are slow



Stationary Wiener filter

$$\hat{d}(n) = \sum_{k=0}^{p} w(k)x(n-k)$$

Solved by the Wiener-Hopf equations

$$R_x w = r_{dx}$$

- (Note: this filter has p+1 coefficients, while in lecture 8 the filter had p coefficients. This is only a difference in notation.)
- x(n) and d(n) were assumed to be jointly WSS



- From the WSS assumption, it follows that the parameters d(n), x(n), x(n-1),..., x(n-p) correlations do not depend on the time n: then the solution neither depends on the time n
- Get rid of the WSS assumption: the solution may change when n changes
- In theory, we can solve the wiener filter for every moment *n*



Let's denote the dependence on n explicitly:

$$\hat{d}(n) = w_n^T x_n$$

where

$$w_n = [w_n(0), w_n(1), \dots, w_n(p)]^T$$
  
 $x_n = [x(n), x(n-1), \dots, x(n-p)]^T$ 

- WH-solution  $w_n = R_x^{-1} r_{dx}$  is the same as previously, but it is separately defined for each *n*
- The solution is not feasible to calculate, because time averages can not be used

- We try to solve the Wiener filter adaptively  $w_{n+1} = w_n + \Delta w_n$ , where  $w_n$  is a solution at the time n
- We require that
  - When the process is WSS, then we must have

$$\lim_{n\to\infty} w_n = R_x^{-1} r_{dx}$$

- $\Delta w_n$  determined only from observations
- We can monitor the changes caused by the nonstationarity



Minimize the error

$$s(n) = E(|e(n)|^2), \ e(n) = d(n) - \hat{d}(n)$$

As in the case of WSS, we get WH equations

$$R_x(n)w_n = r_{dx}(n)$$

- Where all variables depend on the time n
- Autocorrelation matrix  $R_x(n)$  with elements

$$R_x(n)_{ij} = E[x(n-j+1)x^*(n-i+1)]$$



- We solve the correction term  $\Delta w_n$  via error gradient
- Assume that  $w_n$  is a solution at the time n
- We calculate  $w_{n+1}$  by adding one term to  $w_n$ , which reduces the error  $s(n) = \mathrm{E}(|e(n)|^2)$
- The error gradient

$$\nabla s(n) = \left[ \frac{\partial}{\partial w_n(0)} s(n), \quad \frac{\partial}{\partial w_n(1)} s(n), \quad \dots, \quad \frac{\partial}{\partial w_n(p)} s(n) \right]^T$$

shows the direction where the error for which the error is growing the fastest



• We move the vector  $w_n$  in the opposite direction of the gradient, i.e. in the direction where the error decreases most rapidly:

$$w_{n+1} = w_n - \mu \nabla s(n)$$

 The positive step length µ determines how far to go to the negative gradient direction



- Summary of the gradient method:
  - 1. select the initial value of  $w_n$  and step length  $\mu > 0$
  - 2. calculate the gradient  $\nabla s(n)$  using the vector  $w_n$
  - 3. Compute the next vector:

$$w_{n+1} = w_n - \mu \nabla s(n)$$

4. increase *n* by one and go to step 2



 The gradient can be calculated by differentiating with respect to the complex conjugate of w

$$\nabla S(n) = \nabla \operatorname{E}(|d(n) - w_n^T x_n|^2)$$

$$= \operatorname{E}(\nabla[|d(n) - w_n^T x_n|^2])$$

$$= \operatorname{E}(e(n)\nabla[(d(n) - w_n^T x_n)^*])$$

$$= -\operatorname{E}(e(n)x_n^*)$$

· By substitution, we obtain the update rule

$$W_{n+1} = W_n + \mu E(e(n)x_n^*)$$



# Gradient method equilibrium point

Assuming WSS, we obtain

$$-\nabla s(n) = \mathbf{E}(e(n)x_n^*) = r_{dx} - R_x w_n$$

and

$$w_{n+1} = w_n + \mu(r_{dx} - R_x w_n)$$

- If  $w_n = R_x^{-1} r_{dx}$ , i.e. WH-solution, then  $w_{n+1} = w_n$
- So we remain at the solution, if we can get there...



- Subject to certain conditions, the algorithm really reaches the Wiener-Hopf solution:
- If the step length satisfies

$$0 < \mu < \frac{2}{\lambda_{max}}$$

• Where  $\lambda_{max}$  is the largest eigenvalue of the autocorrelation matrix  $R_x$ , then

$$\lim_{n\to\infty} w_n = R_x^{-1} r_{dx}$$



- So, for sufficiently small step size µ, we reach the Wiener-Hopf solution
- The result is true only when x(n) and d(n) are jointly WSS
- For a non stationary process, convergence can not be demonstrated



- Significance of the gradient method is that it can be demonstrated to converge
- Correction term  $\Delta w_n = \mu \operatorname{E}(e(n)x_n^*)$  is however not possible to calculate from the observations in the nonstationary case
- Previously we required that the adaptive filter correction term can be calculated from the observations
- We replace the expectation by its estimate, which is formed as the time average of the values already detected



By selecting L observations, we can estimate

$$E(e(n)x_n^*) = \frac{1}{L} \sum_{l=0}^{L-1} e(n-l)x_{n-l}^*$$

· By substituting in to the gradient method we obtain

$$w_{n+1} = w_n + \frac{\mu}{L} \sum_{l=0}^{L-1} e(n-l) x_{n-l}^*$$



• When L = 1, we get the LMS algorithm:

$$W_{n+1} = W_n + \mu e(n) x_n^*$$

For each component separately we get

$$W_{n+1}(k) = W_n(k) + \mu e(n)x^*(n-k)$$

• Which shows the simplicity of the LMS algorithm. Each iteration requires only the following calculations: for calculating the scalar  $\mu e(n) = \mu(d(n) - w_n^T x_n)$  we need p+1 multiplications, p+1 additions and one multiplication by the constant  $\mu$ .



The difference with the gradient method:

$$\nabla s(n) = -\operatorname{E}(e(n)x_n^*)$$
$$\hat{\nabla}s(n) = -e(n)x_n^*$$

 LMS does not always proceed in the right direction: on the other hand

$$E(\hat{\nabla}s(n)) = -E(e(n)x_n^*) = \nabla s(n)$$

 so on average LMS is progressing in the direction of negative gradient



- Consider how the LMS algorithm progresses for different coefficients µ
- Set the initial values for the parameters as zeros
- The correct solution is w(1) = 1.5 and w(2) = -0.6
- Compare the step lengths  $\mu 1 = 0.002$  and  $\mu 2 = 0.01$
- Demos: Ims1.R and Ims2.R



- Judging by demos, the LMS algorithm converges towards the correct values
- Convergence rate seems to depend on the length of the step
- After convergence, with different step sizes, the LMS algorithm "wobbles" more or less around the true value



- In the LMS algorithm  $w_n$  is a random vector (depending on e(n)and the process x(n)
- Assume that x(n) and d(n) are jointly WSS
- We want to know the conditions under which the LMS algorithm on average converges towards of the Wiener-Hopf solution

$$w = R_x^{-1} r_{dx}$$

i.e. when

$$\lim_{n\to\infty} \mathrm{E}(w_n) = w$$



Taking the expectation of the update rule:

$$E(w_{n+1}) = E(w_n + \mu e(n)x_n^*)$$

$$= E(w_n) + \mu E(d(n)x_n^*) - \mu E\left(x_n^* x_n^T w_n\right)$$

$$= E(w_n) + \mu r_{dx} - \mu E\left(x_n^* x_n^T w_n\right)$$

• Assume that  $w_n$  and  $x_n$  are independent, then

$$E\left(x_n^* x_n^T w_n\right) = R_x E(w_n)$$



We get

$$E(w_{n+1}) = E(w_n) + \mu(r_{dx} - R_x E(w_n))$$

- This is now the gradient method for the vector  $E(w_n)$
- Then we get the LMS algorithm convergence result: with WSS and independence assumptions, the LMS algorithm converges in expectation,  $0<\mu<2/\lambda_{max}$
- Convergence in expectation means that  $E(w_n)$  converges towards the correct value of w



- The convergence result is difficult because the largest eigenvalue  $\lambda_{max}$  of the matrix  $R_x$  should be calculated
- Replace  $\lambda_{max}$  with a larger value, then the coefficient  $\mu$  is certainly between the required values
- Matrix trace is easy to calculate: it can be shown

$$\operatorname{tr}(R_x) = \sum_i \lambda_i \ge \lambda_{max}$$

- eigenvalues therefore not required to be calculated.
- Replace  $\lambda_{max}$  with the trace of  $R_x$



• For a WSS process x(n),  $R_x$  is a Toeplitz matrix and its trace is

then

$$(p+1)r_x(0) = (p+1) E(|x(n)|^2)$$

We get

$$0 < \mu < \frac{2}{(p+1) \operatorname{E}(|x(n)|^2)}$$

• • the variance  $E(|x(n)|^2)$  is much simpler to estimate than  $\lambda_{max}$ 





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- Let's return to the WSS process:
  - 1. expectation is time-independent constant  $m_x = E(x(n))$
  - 2. autocorrelations are independent, i.e. can be written in the form  $r_x(k)$
  - 3. variance is finite:  $c_x(0) < \infty$



- WSS assumptions concern the expectations E(x(n)) and E(x(n))
   x\*(n))
- An important application of modeling is the prediction of future values from the observed values
- Then we need conditional statistics, such as  $E(x(n) \mid x(n-1), x(n-2),...)$
- Conditionality means that the expected value is calculated while assuming that the values of the variables to the right of the vertical line are known



- We previously saw that the conditional expectation of an ARMA process depends on the observed values
- This was used to predict future values
- But for a normally distributed ARMA process, the conditional variance is constant, i.e. the variance of the prediction is always the same regardless of the values observed



AR(1): 
$$x(n) = -ax(n-1) + b(0)v(n)$$
,  $v(n) \sim N(0,1)$ 

Conditional statistics:

$$E(x(n)|x(n-1)) = -ax(n-1)$$

$$\operatorname{var}(x(n)|x(n-1)) = b^2(0)$$

$$\mathrm{E}(x(n))=0$$

$$var(x(n)) = b^{2}(0)/(1 - a^{2})$$

 We see in particular that the conditional variance does not depend on the observations!



- The conditional variance is thus a constant: is this a feature of WSS processes, or only of the normally distributed ARMA processes?
- Example: define a process

$$x(n) \sim \begin{cases} N(0,1), & x(n-1) \ge 0 \\ N(0,2), & x(n-1) < 0 \end{cases}$$



Properties (in the exercises):

$$E(x(n)) = 0$$

$$E(x(n)|x(n-1)) = 0$$

$$var(x(n)) = 1.5$$

$$var(x(n)|x(n-1)) = 0.5 * (-sgn(x(n-1)) + 3)$$

x(n) is WSS



 It is therefore possible that for a WSS process, the conditional variance depends on the observations

• In the previous example, the process is white noise, but its variance can be predicted



- In the ARCH model, the conditional variance is dependent on the observations
- The variance is modeled parametrically on the previous values
- ARCH process is obtained by multiplying the white noise v(n) by the time-dependent standard deviation  $[h(n)]^{1/2}$
- The variance h(n) is defined as a function of the observations  $x(n-1), x(n-2), \ldots$



• ARCH(1)-model 
$$x(n) = v(n)[h(n)]^{1/2}$$

$$h(n) = a(0) + a(1)x^{2}(n-1), \ a(0) > 0, a(1) \ge 0$$

- Noise v(n) values are independent and identically distributed, with expected value of zero and variance 1
- Let's calculate the conditional expectation and variance:

$$E(x(n)|x(n-1)) = 0$$
$$var(x(n)|x(n-1)) = h(n)$$



- The conditional variance is thus directly *h*(*n*) which is not independent of the observations
- Because  $h(n) = a(0) + a(1)x^2(n-1)$  then a high value of  $x^2(n-1)$  causes a large variance at time n, and vice versa
- In financial applications, they talk about volatility clustering (the process gets large values in some time range(s), and small in others)



- Let's see how the ARCH process behaves
- Simulate the above process x(n) with positive coefficients a(0) = 0.03, a(1) = 1.0
- Demo: archex.R, order q = 1
- Demo: archexp.R, order q = 10,
- a(0) = 0.01, a(1) = ... = a(10) = 0.1



More generally, we can define an ARCH(q) process

$$x(n) = v(n)[h(n)]^{1/2}$$

$$h(n) = a(0) + \sum_{k=1}^{q} a(k)x^{2}(n-k)$$

$$a(0) > 0, a(1) \ge 0, \dots, a(q) \ge 0$$

- Noise v(n) defined as in the ARCH(1)-model
- Noise assumption means that the conditional statistics are determined only by h(n)



- Interpretation: If you know the previous findings, the observation x(n) depends only on the value of h(n) calculated from them (and of course the value of the noise v(n))
- The process h(n) is a variance process, because it determines the conditional variance of the ARCH(q) process at time n
- If the parameters a(k) are known, then the variance of h(n) can be calculated from the observations at time n-1
- The ARCH model therefore allows prediction of the variance



- The ARCH model can be generalized in several ways
- The variance function h(n) may depend on the observations in other ways than the ones presented above
- The modeling of the conditional expectation can be included in the model (now the expectation is zero)

