MODEL IDENTIFICATION & DATA ANALYSIS

PART 1

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April 7, 2018

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0 Introduction

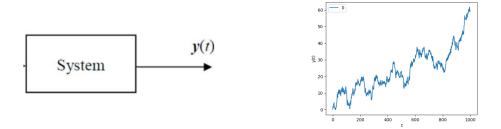
The course will deal with two types of situations:

- 1. Analysis and modelling of **Time-Series**
- 2. Analysis and modelling of Input/Output Systems

0.1 Time-Series

Time series consider vectors $\{y(1), y(2), ..., y(N)\}$ of **measured data** of cardinality N (large , 1000 - 10000).

Said vectors are considered in the **time-domain**: y(t) is a signal or **stochastic process** generated by the system whose output is than sampled.



0.1.1 TS Applications

TS are used for two problems:

- 1. **Prediction problem** : $\{y(1)...y(N)\} \rightarrow \hat{y}(\frac{N+K}{N})$ Given N measurements **estimate** the measurement K timesteps ahead
- 2. Filtering problem : $\{x_1(t)...x_N(t)\} \rightarrow \hat{x}(\frac{t}{t})$ Where $\{x_1(t)...x_N(t)\}$ are internal variables of the system

0.2 I/O Systems

I/O systems consider two measurements :

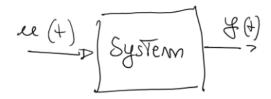
• Input : $\{u(1)...u(N)\}$

• Output: $\{y(1)...y(N)\}$

Resulting in two signals u(t) and y(t). Input signal u(t) can be of two types:

-Controllable : can be affected (ex : voltage)

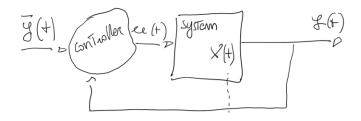
-Uncontrollable : cannot be affected (ex : rain)



0.2.1 I/O Applications

I/O systems are used for three problems :

- 1. **Prediction problem** : $\{y(1)...y(N)\} \rightarrow \hat{y}(\frac{N+K}{N})$ Given N measurements **estimate** the measurement K timesteps ahead
- 2. Filtering problem : $\{x_1(t)...x_N(t)\} \rightarrow \hat{x}(\frac{t}{t})$ Where $\{x_1(t)...x_N(t)\}$ are internal variables of the system
- 3. System control problem: given a desired output $\bar{y}(t)$, control u(t) so that y(t) is as close as possible to $\bar{y}(t)$



0.3 Time Series vs I/O Systems

In prediction and filtering problems both I/O systems and TS can be used. How to chose which one to use?

Ex.1

- **System** : Electric Motor

- **Input** : Current , temperature of motor, electromagnetic fields nearby...

- Output : Torque

We can say that our main input variable (current) is responsible for 90% of the output. The other variables only have slight effects on the torque so they are considered **noise**

The best model to choose is the I/O

Ex.1

- System : Macro-Economic System

- Input : Too many

- Output : Stock prices of FCA

There are many thousand variables affecting the output. Listing and measuring them all would make the model too complex. In this case all the input variables are considered **noise**: the best model to choose is the **Time Series**

Ex.3

- **System**: Environment

- Input : Rain, wind, heatings, cars , temperature, pressure...

- Output : PM10 levels

In this case some main inputs variables can be selected (ex :cars , heating and rain) while the others are modelled as noise. In this case \mathbf{I}/\mathbf{O} model should be used.

It is not wrong to consider all the inputs as noise and model the problem as **Time Series**.

General rule:

		Advantages	
1	TS	Only y(t) must be measured	
	I/O Better estimation		

0.4 Modelling structures

Depending on the problem 2 modelling structures are used.

The TS are modelled with a **mathematical model** which outputs signal y(t). An **imaginary** input e(t) called **white noise** is considered as **standard input** and it is **part of the model**.

The I/O system is modelled by two **mathematical models** which output signal y(t). As above **white noise** is considered as input of one of the two models. The other model has input u(t) which is **not** part of the model.



Figure 1: TS Model

Figure 2: IO Model

All signals and systems are ${\bf time\text{-}discrete}$. Analogue signals are converted to digital signals through ${\bf ADCs}$.

Discrete time points are spaced evenly at pace $\Delta T = \text{sampling time}$

0.5 Mathematical Models

The mathematical models used to elaborate output functions are either white boxes or black - boxes.

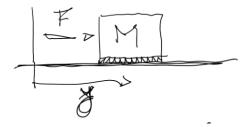


Figure 3: System to be modelled

0.5.1 White Box Models

Also called *first-principles models* assume that the parameters involved in the system are known and well defined. Using white box models we get a **physical interpretation of the model** which makes them useful if the aim is to design the system.

In the example we can derive laws that define our system's **transfer function** given as input a force \vec{F} and output y:

$$\ddot{y}M = F - c\dot{y} \rightarrow Laplace \rightarrow s^2 My = F - scy$$

$$(s^2 M + sc)y = F$$

$$y = \frac{1}{s^2 M + sc}F$$

0.5.2 Black Box Models

In black box models we don't know the internal parameters that influence the system. In our example, we only know that by changing the input \vec{F} a corresponding change in output y(t) can be measured. By measuring the data we can derive a model:

$$y(t) = \frac{b_0 Z^2 + b_1 Z + b_2}{a_0 Z^2 + a_1 Z + a_2} F(t)$$

where $a_0, ..., a_2, b_0, ..., b_2$ are the parameters.

0.5.3 White box vs Black box

Table 1: WB/BB Comparison

White Box	Black Box
-Get physical interpretation of the model and its parameters. -Useful for designing the system	-Very fast -Very accurate -Does not require know-how of the domain -Can be easily re-tuned

0.6 Stochastic Processes

Random variable RV:

v(s) is completely defined by its probability distribution (Gaussian, Uniform...) which is related to its **probability density function** (PDF)

Stochastic Process:

is a sequence of **time-ordered random variables** defined at the same experiment S

where t is the time index. If the experiment is fixed $S = \bar{S}$, we get an instance, a realisation of the stochastic process:

$$v(1,\bar{S}),...,v(t,\bar{S})$$

resulting in a set of samples $\{y(1),...,y(N)\}=\{y(1,\bar{S}),...,y(N,\bar{S})\}$

0.6.1 Characteristics

Mean value m(t):

expected value of a random variable v(t,S) at time t

$$m(t) = E[v(t, S)]$$

Covariance Function $\gamma(t_1, t_2)$:

expected value of the **product** of two **unbiased** random variables at time instants t1, t2:

$$\gamma(t_1, t_2) : E[(v(t_1, S) - m(t_1))(v(t_2, S) - m(t_2))]$$

Removing the mean brings the signal closer to 0.

If t1 = t2 = t the covariance degenerates in **variance**:

$$\gamma(t) = E[(v(t, S) - m(t))^2]$$

0.6.2 Stationary Stochastic Processes

Has properties:

- 1. $m(t) = m, \forall t$
- 2. $\gamma(t_1, t_2)$ depends on $\tau = |t_1 t_2|$

This means that the covariance depends on the **distance in time** and not on specific considered samples.

$$\gamma(t_1, t_2) = \gamma(t_3, t_4) \to |t_1 - t_2| = |t_3 - t_4|$$

 $\gamma(\tau) = E[(v(t) - m)(v(t - \tau) - m)]$ has properties :

- $\gamma(0) = E[(v(t) m)^2] \rightarrow$ variance
- $\bullet \ |\gamma(\tau)| \leq \gamma(0)$
- $\gamma(\tau) = \gamma(-\tau)$

SSP Equivalence

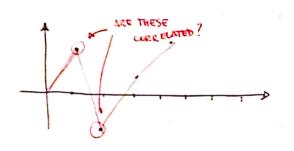
Two SSPs $y_1(t), y_2(t)$ are equivalent in a **weak sense** if:

- $m_{v1} = m_{v2}$
- $\gamma_{y1}(\tau) = \gamma_{y2}(\tau) , \forall \tau$

Correlation Function

If the m=0 the $\gamma(\tau)$ function degenerates in the **correlation function**:

$$E[v(t)v(t-\tau)]$$



0.6.3 White Noise

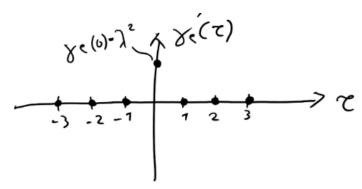
e(t) is SSP called white noise and is written as

$$e(t) \to WN(\mu, \lambda^2)$$

Properties:

• Mean value : $E[e(t)] = \mu, \forall t$

• Variance : $\gamma_e(0) = E[(e(t) - \mu)^2] = \lambda^2$



No covariance means that the samples are **not related** Considering a **Gaussian Distribution** : $e(t) \to WGN(\mu, \lambda^2)$

0.7 Sample estimation of mean and covariance function

Dealing with samples it is useful to **estimate** the mean and covariance of the samples.

Output y(t) is a SSP : $\{y(1), ..., y(N)\}$ a particular realisation of \bar{S} with :

- Mean m = E[y(t)]
- Covariance $\gamma(\tau) = E[(y(t) m)(y(t \tau) m)]$

This seems trivial but the computation of the expected value cannot be done because the **distribution of the process** is **unkown**.

These two can be estimated

0.7.1 Sample Mean

The sample mean is a good estimator for the mean m:

$$\hat{m}_n = \frac{1}{N} \sum_{t=1}^{N} y(t)$$

Properties of the estimator:

1. \hat{m}_n is **correct** if $E[\hat{m}_n] = m$

Proof:
$$E[\hat{m}_n] = E[\frac{1}{N} \sum_{t=1}^{N} y(t,s)] = \frac{1}{N} \sum_{t=1}^{N} E[y(t,s)] = \frac{1}{N} \sum_{t=1}^{N} m = \mathbf{m}$$

Example

$$y(t,S) = \bar{v}(s) \to WN(0,1) \text{ and } S = \bar{S}, \{y(1,\bar{S}), ..., y(N,\bar{S})\} \text{ so }:$$

$$-\hat{m}_n = \frac{1}{N} \sum_{t=1}^N y(t,\bar{S}) = \frac{1}{N} \sum_{t=1}^N \bar{v}(\bar{S}) = \frac{1}{N} N \bar{v}(\bar{S}) \neq 0 \to \text{ bad estimator}$$

$$-\check{m}_n = \frac{1}{N} \sum_{S=1}^N y(\bar{t},S) = \frac{1}{N} v(S) \to 0 \to \text{ good estimator}$$

2. \hat{m}_n is **consistent** if $E[(\hat{m}_n - m)^2] \xrightarrow[N \to \infty]{N \to \infty} 0$

The **error variance** approaches 0 for large values of N: this means that with a lot of data $N \to \infty$ we can estimate \hat{m}_n more effectively.

In general one can say that \hat{m}_n is consistent if $\gamma(\tau) \xrightarrow[|\tau| \to \infty]{} 0$

Example

$$y(t,S) = \bar{V}(S) \to WN(0,1)$$

 $\gamma(\tau) = E[(\gamma(\tau))(\gamma(t-\tau))] = E[\bar{V}(S)\bar{V}(S)] = E[\bar{V}(S)^2] = 1$

0.7.2 Sample Covariance

y(t) is a SSP with **zero mean**.

A good estimator for for the covariance is the **sample covariance**:

$$\hat{\gamma}_N(\tau) = \frac{1}{N-\tau} \sum_{t=1}^{N-\tau} y(t)y(t+\tau)$$
$$0 \le \tau \le N-1$$

It is important to notice that this approximation is good for $\tau \ll N$ because the accuracy of $\gamma_N(\tau)$ decreases with τ

Properties of the estimator:

1. $\hat{\gamma}_N(\tau)$ is **correct** if $E[\hat{\gamma}_N(\tau)] = \gamma(\tau)$

Proof:

$$E[\hat{\gamma}_N(\tau)] = E\left[\frac{1}{N-\tau} \sum_{t=1}^{N-\tau} y(t)y(t+\tau)\right] = \frac{1}{N-\tau} \sum_{t=1}^{N-\tau} E[y(t)y(t+\tau)] = \frac{1}{N-\tau} \sum_{t=1}^{N-\tau} \gamma(\tau) = \gamma(\tau)$$

2.
$$\hat{\gamma}_N(\tau)$$
 is **consistent** if $E[(\hat{\gamma}_N(\tau) - \gamma(\tau))^2] \xrightarrow[N \to \infty]{} 0$, True if $\gamma(\tau) \xrightarrow[|\tau| \to \infty]{} 0$

Observation 1:

$$\hat{\gamma}_N(\tau) = \frac{1}{N-\tau} \sum_{t=1}^{N-\tau} y(t)y(t+\tau)$$

$$0 \le \tau \le N-1, \tau \ge 0$$
but since y(t) is a SSP $\gamma(\tau) = \gamma(-\tau)$:

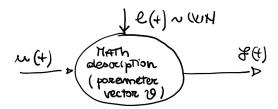
$$\hat{\gamma}_N(\tau) = \frac{1}{N - |\tau|} \sum_{t=1}^{N - |\tau|} y(t)y(t + |\tau|)$$
$$|\tau| \le N - 1$$

Observation 2:
$$\hat{\gamma}'_N(\tau) = \frac{1}{N-|\tau|} \sum_{t=1}^{N-|\tau|} y(t) y(t+|\tau|) \to E[\hat{\gamma}'_N(\tau)] = \dots = \frac{1}{N} \gamma(\tau) (N-|\tau|)$$
 As shown $\hat{\gamma}'_N(\tau)$ doest not satisfy the correct property.

However for N $\to \infty$ and $\tau << N$: $\hat{\gamma}'_N(\tau)$ is asimptotically correct

1 Chapter 1

1.1 Model classes



$$\mbox{Mathematical model} = \begin{cases} u(t) & \mbox{input (I/O only)} \\ e(t) & \mbox{white noise} \\ y(t) & \mbox{output} \end{cases}$$

The mathematical model is described by **parametric parameter vector** θ that is found using a **parametric supervised** identification approach.

The models can be described with:

- Differential Equations in time domain
- Transfer functions

1.1.1 Time-Series model classes

The following processes are modelled with differential equations

1. Moving Average Models (MA):

A process y(t) **generated** by a WN e(t) is a moving average of order n MA(n) process if:

$$y(t) = c_0 e(t) + c_1 e(t-1) + \dots + c_n e(t-n)$$

with parameter vector $\theta = \{c_0, ..., c_n\}$

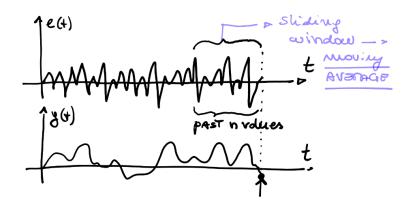


Figure 4: y(t) is linear combination of past n e(t) values

2. Autoregressive Models (AR):

A process y(t) **generated** by a WN e(t) is an autoregressive of order m AR(m) process if:

$$y(t) = a_1 y(t-1) + a_2 y(t-2) + \dots + a_m y(t-m) + c_0 e(t)$$

with parameter vector $\theta = \{c_0, a_1, ..., a_m\}$

3. Autoregressive Moving Average Models (ARMA):

A process y(t) **generated** by a WN e(t) is an ARMA of order (n,m) **ARMA(n,m)** process if:

$$y(t) = a_1 y(t-1) + a_2 y(t-2) + \dots + a_m y(t-m) + c_0 e(t) + c_1 e(t-1) + \dots + c_n e(t-n)$$

with parameter vector $\theta = \{c_0, ..., c_n, a_1, ..., a_m\}$

 $ARMA(0,n) \rightarrow MA(n) : MA(n) \text{ is } \mathbf{subclass} \text{ of } ARMA$

 $ARMA(m,0) \rightarrow AR(m) : AR(m) \text{ is } \mathbf{subclass} \text{ of } ARMA$

1.1.2 Input/Output model classes

The following processes are modelled with differential equations

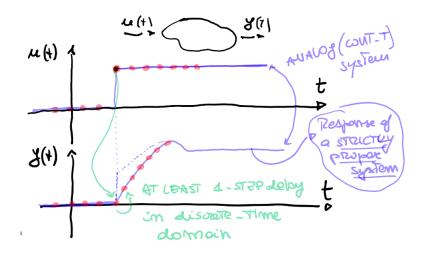
1. Autoregressive Moving Average Exogenous (ARMAX):

A process y(t) **generated** by a WN e(t) and **exogenous** signal u(t) is an ARMAX of order (n,m, p+k) process if:

$$y(t) = a_1 y(t-1) + \dots + a_m y(t-m) + c_0 e(t) + \dots + c_n e(t-n) + b_0 u(t-k) + \dots + b_p u(t-k-p)$$

with parameter vector $\boldsymbol{\theta} = \{c_0, ..., c_n, a_1, ..., a_m, b_0, ..., b_p\}$

 $K \geq 1$ plays an important role : it represents the pure/intrinsic **delay** between y(t) and u(t). If u(t) is a step the corresponding output y(t) is



shown in figure.

Sampling (red dots) gives a discrete approximation: when the input slope rises a sample is taken resulting in a high value. The corresponding output is still low: this causes a **1** step delay

Example

$$y(t) = \frac{1}{2}y(t-1) + \frac{1}{3}y(t-2) + e(t) + e(t-3) + u(t-2) + \frac{1}{2}u(t-4)$$
 The process is an ARMAX $(2,3,2+2)$

Observation: missing values as above can be present!

Remark

Armax models are the most general class models for **dynamic**, **linear**, **time-invariant** systems.

Non-Linear N-ARMAX y(t) = f(y(t-1), ..., y(t-m), e(t), ..., e(t-n), u(t-k), ..., u(t-k-p)

depend on **non-linear functions** : polynomials , splines ,NN ,Radial Basis Functions ,Fuzzy Sets.

1.2 Transfer function representation

The four models found above can be represented using **transfer functions**. To transform time domain equations into the equivalent transfer function representation the **Z** operator is introduced.

1.2.1 Z Operator

• The operator Z^{-1} is the **backward shift** operator :

$$Z^{-1}x(t) = x(t-1)$$

• The operator Z^{+1} is the **forward shift** operator :

$$Z^{+1}x(t) = x(t+1)$$

Both operators have properties:

• Linearity : $Z^{-1}(ax(t) + by(t)) = Z^{-1}ax(t) + Z^{-1}by(t) = ax(t-1) + by(t-1)$

• Recursion : $Z^{-1}(Z^{-1}...(Z^{-1}x(t))) = x(t-n) = Z^{-n}$

1.2.2 Time domain to Transfer Function

The Z operators are used to shift the equations of the time domain to be all at time \mathbf{t} .

In case of a generic ARMAX(m,n,p+k) process

$$y(t) = a_1 y(t-1) + \ldots + a_m y(t-m) + c_0 e(t) + \ldots + c_n e(t-n) + b_0 u(t-k) + \ldots + b_p u(t-k-p)$$

Applying the Z^{-1} operator:

$$y(t) = a_1 Z^{-1} y(t) + \dots + a_m Z^{-m} y(t) + c_0 e(t) + \dots + c_n Z^{-n} e(t) + b_0 Z^{-k} u(t) + \dots + b_p Z^{-k-p} u(t)$$

Collecting:

$$y(t)[1 - a_1Z^{-1} + \dots + a_mZ^{-m}] = [c_0e + \dots + c_nZ^{-n}]e(t) + [b_0Z^{-k} + \dots + b_nZ^{-k-p}]u(t)$$

Dividing:

$$y(t) = \frac{[c_0e + \dots + c_nZ^{-n}]}{[1 - a_1Z^{-1} + \dots + a_mZ^{-m}]}e(t) + \frac{[b_0 + \dots + b_pZ^{-p}]}{[1 - a_1Z^{-1} + \dots + a_mZ^{-m}]}u(t)Z^{-k}$$

Defining:

$$A(Z) = 1 - a_1 Z^{-1} + \dots + a_m Z^{-m}$$

$$B(Z) = b_0 + \dots + b_p Z^{-p}$$

$$C(Z) = c_0 e + \dots + c_n Z^{-n}$$

The resulting process using TF representation is:

$$y(t) = \frac{C(Z)}{A(Z)}e(t) + \frac{B(Z)}{A(Z)}u(t)Z^{-k}$$



1.2.3 From Z^- to Z^+

The transfer functions can be written in negative, positive o mixed power of Z. The example explains how to get the positive power representation starting from a negative one:

$$y(t) = \frac{c_0 + c_1 Z^{-1} + \dots + c_n Z^n}{1 - a_1 Z^{-1} - \dots - a_m Z^{-m}} e(t)$$

If $m \ge n$ by multiplying by Z^{+m} :

$$y(t) = \frac{c_0 Z^m + c_1 Z^{m-1} + \dots + c_n Z^{m-n}}{Z^m - a_1 Z^{m-1} - \dots - a_m} e(t)$$

Observation

Even if feasible and correct it is better to avoid the mixed representation!

1.2.4 Importance of stationary property

Transformation **Time Domain** \leftrightarrow **Transfer Functions** are **feasible** if the **stationary property** holds because otherwise the Z operator is not applicable.

$$y(t) = \frac{Z + \frac{1}{2}}{Z - \frac{1}{2}}e(t), e(t) - WN(0,1)$$

In time domain

$$(Z - \frac{1}{3})y(t) = (Z + \frac{1}{2})e(t)$$
$$y(t+1) - \frac{1}{3}y(t) = e(t+1) + \frac{1}{2}e(t)$$
$$y(t+1) = \frac{1}{3}y(t) + e(t+1) + \frac{1}{2}e(t)$$

Time shift to start a time "t" (can be done in stationary conditions):

$$y(t) = \frac{1}{3}y(t-1) + e(t) + \frac{1}{2}e(t-1)$$

Back to TF

$$y(t) = \frac{1}{3}Z^{-1}y(t) + e(t) + \frac{1}{2}Z^{-1}e(t)$$
$$[1 - \frac{1}{3}Z^{-1}]y(t) = [1 + \frac{1}{2}Z^{-1}]e(t)$$
$$y(t) = \frac{[1 + \frac{1}{2}Z^{-1}]}{[1 - \frac{1}{3}Z^{-1}]}e(t)$$

1.2.5 Pole, Zeros and Stability

Considering a process with signals e(t), y(t) and a system W(Z) represented in **positive/null** power:

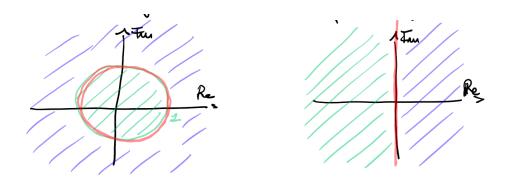
- Poles of W(Z) are the roots of the denominator
- **Zeros** of W(Z) are the **roots** of the nominator

A system is said to be **asymptotically stable** if and only if all the **poles** of W(Z) are **strictly inside** the unit circle (left graph).

Blue = unstable region

Red = simple stability region

Green = asymptotically stability region



Note: if we were dealing with **continuous** signals and processes instead of Z transformation we would apply **Laplace**. Also the stability region would change as seen on the right graph.

Example:

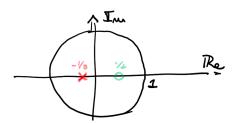
$$W(Z) = \frac{1 - \frac{1}{2}Z^{-1}}{1 + \frac{1}{3}Z^{-1}}$$

Move to positive power:

$$W(Z) = \frac{Z - \frac{1}{2}}{Z + \frac{1}{3}}$$

• Pole : $Z=-\frac{1}{3}$

• Zero : $Z = \frac{1}{2}$



The system is asymptotically stable since all poles are within the unit circle.

1.2.6 Stationary property and stability

In a stochastic process y(t) obtained as output of a system W(Z) fed with a stochastic process v(t), y(t) is a **stationary process SSP** if and only if:

- 1. v(t) is a stationary stochastic process
- 2. W(Z) is asymptotically stable

Checking the stationary property is usually very long, instead these two properties make it easy: input v(t) is usually a **white noise** which is a **stationary stochastic process**.

1.2.7 Poles and Zeros in MA & AR processes

MA(1)

$$y(t) = e(t) + \frac{1}{2}e(t-1)$$
$$y(t) = (1 + \frac{1}{2}Z^{-1})e(t)$$
$$y(t) = (\frac{Z + \frac{1}{2}}{Z})e(t)$$

- **Zero** : $Z = -\frac{1}{2}$
- **Pole** : Z = 0

AR(1)

$$y(t) = \frac{1}{2}y(t-1) + 3e(t)$$
$$y(t) = (\frac{3Z}{Z - \frac{1}{2}})e(t)$$

- **Zero** : Z = 0
- **Pole** : $Z = \frac{1}{2}$

General conclusion

A MA(n) process is generated by a TF having:

- n generic zeros
- n poles all in $0 \rightarrow$ always stationary!

It is also called **All-Zeros** process

An **AR(m)** process is generated by a TF having:

- m zero all in 0
- m generic poles It is also called All-Poles process

2 Chapter 2: Analysis of Stochastic Processes

TS modelled with ARMA processes and I/O modelled with ARMAX models can be represented in 4 different ways :

- Time domain (Chap.1)
- Transfer function (Chap.1)
- Probabilistic representation
- Frequency representation

2.1 Probabilistic Representation

2.1.1 Probabilistic representation of MA(n)

Time domain representation : $y(t) = c_0 e(t) + ... + c_n e(t-n), e(t) \sim WN(0, \lambda^2)$. The process is **stationary** as all the poles are in the origin.

• Mean of y

$$m_y = E[y(t)] = E[c_o e(t) + \dots + c_n e(t-n)] = c_0 E[e(t)] + \dots + c_n E[e(t-n)]$$

Because of stationary property $E[e(t)] = \dots = E[e(t-n)] = 0$

$$m_y = 0$$

• Covariance of y

$$-\tau = 0$$

$$\gamma_y(0) = E[(y(t) - m_y)^2] = E[y(t)^2] = E[(c_0e(t) + \dots + c_ne(t-n))^2] = c_0^2 E[e(t)^2] + \dots + c_n^2 E[e(t-n)^2] + 2c_0c_1 E[e(t)e(t-1) + \dots + 2c_{n-1}c_n E[e(t-n-1)e(t-n)]$$
where
$$E[e(t)^2] = E[e(t-1)^2] = \dots = E[e(t-n)^2] = \lambda^2$$

$$E[e(t)e(t-1)] \dots = 0 \text{ because not correlated}$$

$$\gamma_y(0) = \lambda^2(c_0^2 + \dots + c_n^2)$$

$$\begin{aligned} & -\tau = 1 \\ & \gamma_y(1) = E[(y(t) - m_y)(y(t-1) - m_y)] = E[y(t)y(t-1)] \\ & E[(c_0e(t) + \ldots + e_ne(t-n))(c_0e(t-1) + \ldots + c_ne(t-n-1)] \\ & \text{only terms at same time survive :} \\ & c_0c_1E[e(t-1)^2] + \ldots + c_{n-1}c_nE[e(t-n)^2] \\ & \text{where } E[e(t-i)^2] = \lambda^2 \end{aligned}$$

$$\gamma_y(1) = (c_0c_1 + c_1c_2 + \dots + c_{n-1}c_n)\lambda^2$$

 $-\tau=2$

$$\gamma_y(2) = (c_0c_2 + c_1c_3 + \dots + c_{n-2}c_n)\lambda^2$$

- ..

 $-\tau = n$

$$\gamma_y(n) = c_0 c_n \lambda^2$$

 $-|\tau|>n$

$$\gamma_y(\tau) = 0, \tau > n$$

Which means that MA(n) has a finite memory of n steps

2.1.2 Probabilistic representation of AR(1)

$$y(t) = ay(t-1) + e(t), e(t) \sim WN(0, \lambda^2)$$

Is $y(t)$ a **SSP**?

TF representation:

$$y(t)=aZ^{-1}y(t)+e(t)\to y(t)=\frac{1}{1-aZ^{-1}}e(t)\to y(t)=\frac{Z}{Z-a}e(t)$$

So $y(t)$ is a SSP if and only if $|a|<1$

• Mean of y

$$m_y = E[y(t)] = E[ay(t-1) + e(t)] = am_y + m_e$$

 $m_y(1-a) = m_e \to m_y = \frac{m_e}{1-a}$

 $m_e = 0$ so:

$$m_y = 0$$

This hold only if the general input v(t) has $m_v = 0$ and the system $\mathbf{W}(\mathbf{Z})$ is asymptotically stable.

• Covariance of y

$$\begin{split} &-\tau = 0 \\ &\gamma_y(0) = E[(y(t) - m_y)^2] = E[y(t)^2] = E[(ay(t-1) + e(t))^2] \\ &\gamma_y(0) = a^2 E[y(t-1)^2] + E[e(t)] + 2aE[e(t)y(t-1)] \\ &\gamma_y(0) = a^2 \gamma_y(0) + \lambda^2 + 0 \end{split}$$

Observation: the fact that E[e(t)y(t-1)] = 0 is explained in 2.1.3!

$$\gamma_y(0) = \frac{\lambda^2}{1 - a^2}$$

$$- \tau = 1$$

$$\gamma_y(1) = E[(y(t) - m_y)(y(t-1) - m_y)] = E[(ay(t-1) + e(t))y(t-1)]$$

$$\gamma_y(1) = E[ay(t-1)y(t-1)] + E[e(t)y(t-1)] = a\gamma_y(0) + 0$$

Observation: the fact that E[e(t)y(t-1)] = 0 is explained in 2.1.3!

$$\gamma_y(1) = a\gamma_y(0)$$

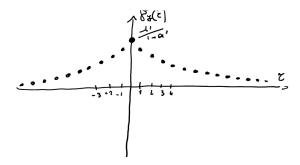
– ...

$$-\tau \neq 0$$

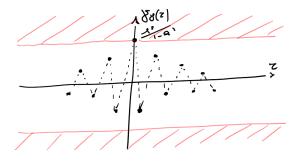
$$a\gamma_y(\tau-1)$$

Which means that AR(1) has an **infinite memory**. The formula is also known as Yule-Walker formula of order 1

1.Plot of $\gamma_y(\tau), 0 < a < 1$



2.Plot of $\gamma_y(\tau), -1 < a < 0$



2.1.3 AR/ARMA as $MA(\infty)$

A general rule states that every AR/ARMA stationary stochastic process can be modelled as $MA(\infty)$. Example with AR(1) as above:

$$y(t) = \frac{1}{1 - aZ^{-1}} e(t) \to y(t) = \sum_{k=0}^{\infty} (aZ^{-1})^k e(t)$$

A geometric series of common ratio aZ^{-1}

$$y(t) = e(t)[1 + aZ^{-1} + a^2Z^{-2} + \ldots]$$

$$y(t) = e(t) + ae(t-1) + a^2e(t-2)...$$

Which is the $\mathbf{MA}(\infty)$ equivalent of AR(1) This formula is very useful to demonstrate that in an AR(1) E[e(t)y(t-1)] = 0 by expressing y(t-1) in $MA(\infty)$: $E[e(t)(e(t-1) + ae(t-2) + a^2e(t-2)...) = E[e(t)e(t-1)] + E[e(t)ae(t-2)... = 0$ Due to correlation all terms are equal to zero (WN property!).

2.2 Frequency Representation

The power density / spectrum of a SSP y(t):

$$\Gamma_y(w) = \sum_{\tau = -\infty}^{\infty} \gamma_y(\tau) e^{-jw\tau}$$

where $\Gamma_y(w)$ is the **Discrete Fourier Transform**.

Properties:

- 1. $\Gamma_y(w)$ is a **real** function of a **real** variable w which means that $Im\{\Gamma_y(w)\}=0$
- 2. $\Gamma_y(w)$ is a **positive** function which means that $\Gamma_y(w) \geq 0, \forall w \in \Re$
- 3. $\Gamma_y(w)$ is an **even** function which means that $\Gamma_y(w) = \Gamma_y(-w)$
- 4. $\Gamma_y(w)$ is a **periodic** function of period 2π which means that $\Gamma_y(w) = \Gamma_y(w + k 2\pi)$.

2.3 Inverse Fourier Transform

Fourier Transform:

$$\Gamma_y(w) = F\{\gamma_y(\tau)\} = \sum_{t=-\infty}^{\infty} \gamma_y(\tau)e^{-jw\tau}$$

Inverse Fourier Transform:

$$\gamma_y(\tau) = F^{-1}\{\Gamma_y(w)\} = \frac{1}{2\pi} \int_{-\pi}^{\pi} \Gamma_y(w) e^{jw\tau} dw$$

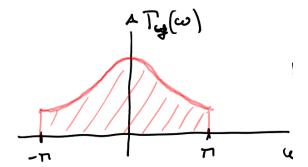
It is important to notice that $\Gamma_y(w)$ and $\gamma_y(\tau)$ contain the same information: passing from one to another does not result in loss or gain of information.

Special IFT: Computation of variance

A special case of IFT is the computation of the variance, when $\tau = 0$:

$$\gamma_y(0) = \frac{1}{2\pi} \int_{-\pi}^{\pi} \Gamma_y(w) dw$$

which is the **area below the spectrum** between $(-\pi, \pi)$ divided by 2π .



2.4 White Noise in the frequency domain

In case we are dealing with a WN : $e(t) = WN(0, \lambda^2)$ we can consider it in three different domains.

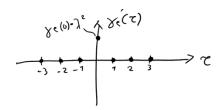
1. Time domain

WN is clearly unpredictable



2. Probabilistic domain

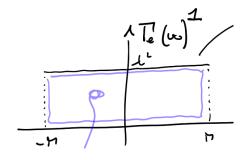
Considering the WN in the probabilistic domain and plotting its **variance** only $\gamma_e(0) \neq 0$: there is no **correlation** between e(t) and $e(t \pm \tau)$



3. Frequency domain

Since the definition of FT relies on the definition of **covariance** $\gamma_e(\tau)$, as seen in point 2 only for $\tau = 0 \rightarrow \gamma_e(\tau) \neq 0$:

$$\Gamma_e(w) = \gamma_e(0)e^{jw0} = \gamma_e(0) = \lambda^2$$



The area is $2\pi\lambda^2$ so the variance is $\frac{area}{2\pi} = \lambda^2 = \gamma_e(0)$ The **energy** of the WN is **uniformly distributed** over all frequencies.

2.5 Computation of the spectrum of a process generated as the output of a digital system

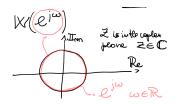
Problem: the **computation** of the $\Gamma_y(w)$ is quite difficult most of the times. A **simpler** alternative can be found using the notion of **Frequency Response**.

2.5.1 Frequency Response of a linear system

Given two signals (SSP) input v(t) and output y(t), where input passes through W(Z) the system or digital filter, then the **frequency response** is

$$W(e^{jw})$$

which corresponds to the evaluation of the **transfer function** on the **unit circumference**



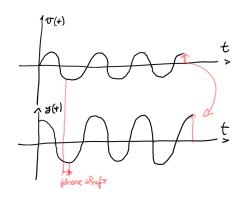
The frequency response is used in system theory in the $\mathbf{Frequency}$ $\mathbf{Response}$ $\mathbf{Theorem}$

FR Th.

If W(Z) is **asymptotically stable** and v(t) is $Asen(\Omega t + \phi)$, where A is the amplitude an ϕ the phase of the sinusoid, the the output is a **pure sinusoid** with:

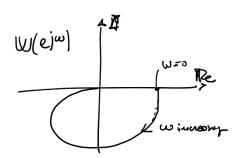
- \bullet the same angular speed Ω
- amplitude $A|W(e^{j\Omega})| \to \mathbf{gain}$
- phase $\phi + \angle W(e^{j\Omega}) \to \mathbf{shift}$ in phase

$$y(t) = A|W(e^{j\Omega})|sen(\Omega t + \phi + \angle W(e^{j\Omega}))|$$



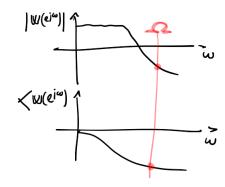
1.FR Nyquist plot

 $W(e^{jw})$ is a complex function of a **real** variable.



2.FR Bode plot

Bode plot gives information about magnitude and phase



2.5.2 Spectrum computation with FR

If y(t) is output of a transfer function W(Z) which is **asymptotically stable** with input signal v(t), then the spectrum is:

$$\Gamma_y(w) = |W(e^{jw})|^2 \Gamma_v(w)$$

The computation of $\Gamma_v(w)$ still remains but most of the time signal v(t) is a white noise $\sim (0, \lambda^2)$ which means that $\Gamma_v(w) = \lambda^2$

2.6 Equivalent representations of ARMA

An ARMA SSP can be represented in 4 different but **equivalent** ways with $e(t) \sim WN(0,1)$:

- 1. Time domain $y(t) = a_1 y(t-1) + ... + a_m y(t-m) + c_0 e(t) + ... + c_n e(t-n)$
- 2. Transfer function $y(t) = \frac{C(Z)}{A(Z)}e(t)$
- 3. Probabilistic domain:

$$\begin{cases} m_y &= E[y(t)] \\ \gamma_y(\tau) &= E[(y(t) - m_y)(y(t - \tau) - m_y)] \end{cases}$$

4. Frequency domain

$$\begin{cases} m_y & = E[y(t)] \\ \Gamma_y(w) & w \in \Re \end{cases}$$

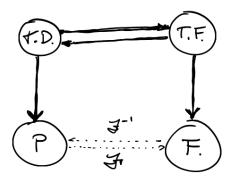
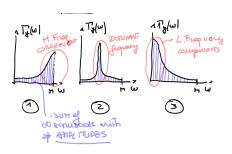


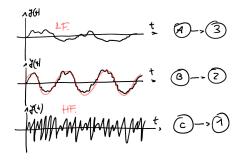
Figure 5: Bold : usual transformation , dotted : feasible but difficult

2.7 Example & Exercises

Example 1

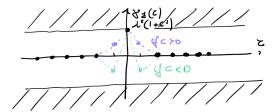
Given 3 output spectra match the corresponding time domain representation.





Example 2

Given a **MA(1)** process $y(t) = e(t) + ce(t-1), e \sim WN(0, \lambda^2), c \in \Re$. - y(t) is **stationary** because MA(1) is always **asymptotically stable** - $m_e = 0 \rightarrow m_y = 0$ - $\gamma_y(0) = \lambda^2(1 + c^2)$ - $\gamma_y(1) = \lambda^2 c$ - $\gamma_y(\tau) = 0, |\tau| \ge 2$



1. Composition of $\Gamma_y(w)$ with $\lambda^2 = 1$:

• From definition

$$\Gamma_y(w) = \sum_{\tau = -\infty}^{\infty} \gamma_y(\tau) e^{-jw\tau}$$

- For $\tau = 0 : (1 + c^2)$
- For $\tau = 1 : ce^{-jw}$
- For $\tau = -1 : c + e^{+jw}$
- For $|\tau| \ge 2:0$

$$\Gamma_{\nu}(w) = 1 + c^2 + c(e^{-jw} + e^{+jw})$$

 $\text{Recall}:\, e^{-jw}+e^{jw}=cosw-jsenw+cosw+jsenw=2cosw$

$$\Gamma_y(w) = (1 + c^2) + 2ccosw$$

Which is real, positive, even, periodic

• From frequency response

The MA(1) transfer function is

$$y(t) = (1 + cZ^{-1})e(t)$$

$$\Gamma_{\nu}(w) = |W(e^{jw})|^2 \Gamma_{e}(w) = |1 + ce^{-jw}|^2 \cdot 1$$

Recall :
$$|a+jb|^2 = Im[a+ib]^2 + Re[a+ib]^2 = a^2 + b^2 = (a+jb)(a-jb)$$

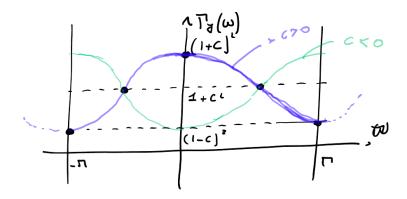
$$(1+ce^{-jw})(1-ce^{jw}) = 1+c^2(e^{jw} \cdot e^{-jw}) + c(e^{-jw} + e^{jw}) = 1+c^2 + 2c\cos w$$

2. Plotting of $\Gamma_y(w)$:

$$\Gamma_y(0) = (1+c)^2$$

$$\Gamma_y(\frac{\pi}{2} = 1 + c^2)$$

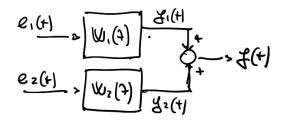
$$\Gamma_{\nu}(\pi) = (1 - c)^2$$



3. Compute the variance $\gamma_y(0)$ given $\Gamma_y(w)$:

$$\gamma_y(0) = \frac{1}{2\pi} \int_{-\pi}^{\pi} \Gamma_y(w) dw = \frac{1}{2\pi} \int_{-\pi}^{\pi} (1 + c^2 + 2ccosw) dw$$
$$\frac{1}{2\pi} \int_{-\pi}^{\pi} (1 + c^2) dw + \frac{1}{2\pi} \int_{-\pi}^{\pi} 2ccosw dw$$
$$\frac{1}{2\pi} [(1 + c^2)[w]_{-\pi}^{\pi} + 2c[senw]_{-\pi}^{\pi}] = \frac{1}{2\pi} [(1 + c^2)(2\pi)] = 1 + c^2$$

Example 3



Consider the SSP y(t) generated by 2 inputs.

 $W_1(t), W_2(t)$ are asymptotically stable.

$$e_1(t) \sim W(0, \lambda_1^2), e_2(t) \sim W(0, \lambda_2^2)$$

$$e_1 \perp e_2 \to E[e_1(t)e_2(t-\tau)] = 0$$

Calculate $\gamma_y(\tau)$ and $\Gamma_y(w)$

$$\begin{split} \bullet \ \, \gamma_y(\tau) \\ \gamma_y(\tau) &= E[y(t)y(t-\tau)] = E[(y_1(t)+y_2(t))(y_1(t-\tau)+y_2(t-\tau))] \\ E[y_1(t)y_1(t-\tau)] + E[y_2(t)y_2(t-\tau)] + E[y_1(t)y_2(t-\tau)] + E[y_2(t)y_1(t-\tau)] \\ \gamma_{y_1}(\tau) + \gamma_{y_2}(\tau) + 0 + 0 \end{split}$$

Term 3 and 4 are = 0 which is a result obtained by rewriting them as $MA(\infty)$ and exploiting the hypothesis that $e_1(t) \perp e_2(t)$.

$$\gamma_y(t) = \gamma_{y_1}(t) + \gamma_{y_2}(t)$$

•
$$\Gamma_y(t)$$

$$\Gamma_y(t) = \sum_{\tau = -\infty}^{\infty} \gamma_y(\tau) e^{-jw\tau} = \sum_{\tau = -\infty}^{\infty} \gamma_{y_1}(\tau) e^{-jw\tau} + \sum_{\tau = -\infty}^{\infty} \gamma_{y_2}(\tau) e^{-jw\tau}$$

$$\Gamma_y(w) = \Gamma_{y_1}(w) + \Gamma_{y_2}(w)$$

The result can be generalised to more than 2 inputs that are summed to form an SSP y(t):

$$[\gamma_y(t) = \gamma_{y_1}(t) + \gamma_{y_2}(t) + \dots + \gamma_{y_k}(t)]$$
$$[\Gamma_y(w) = \Gamma_{y_1}(w) + \Gamma_{y_2}(w) + \dots + \Gamma_{y_k}(t)]$$

The result hold if all $W_i(t)$ are asymptotically stable, all $v_i(t)$ are ssp and incorrelated

Example 4

Consider the following AR(1) SSP $y(t) = \frac{1}{3}y(t-1) + e(t) + 2 \rightarrow e \sim WN(1,1)$ which has a asymptotically stable TF.

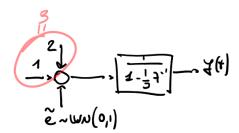
Calculate m_y and γ_y .

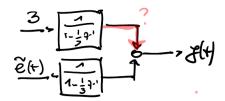
• Mean of y

- Method 1
$$E[y(t)] = E[\frac{1}{3}y(t-1)e(t)+2] \to (1-\frac{1}{3})m_y = m_e+2 \to m_y = \frac{9}{2}$$
 - Mehotd 2

$$e(t) = \tilde{e} + 1, \tilde{e} \sim WN(0, 1)$$

Using the superposition principle of LTI systems:





The constant value 3 can be seen as a sinusoidal signal with **w=0** so the Frequency Response Theorem can be applied:

$$3(\frac{1}{1-\frac{1}{3}Z^{-1}})$$
 calculated in $z=e^{j0}$

$$3(\frac{1}{1-\frac{1}{3}}) = \frac{9}{2}$$
 And since $m_{\tilde{e}}$ is a zero mean signal $\rightarrow m_y = \frac{9}{2}$

• Covariance of y

- Method 1: BAD

$$\begin{split} E[(y(t)-\tfrac{9}{2})^2] &= E[(\tfrac{1}{3}y(t-1)+e(t)+2-\tfrac{9}{2})^2] \\ \gamma_y(0) &= \tfrac{1}{9}E[y(t-1)^2] + E[e(t)^2] + \tfrac{25}{4} + \tfrac{2}{3}E[y(t-1)e(t)] - \tfrac{5}{3}E[y(t-1)] + 5E[e(t)] \end{split}$$

Remark:
$$E[(e(t) - m_e) = \gamma_e(0) = E[e(t)^2] - 2E[e(t)m_e] + m_e^2$$

 $E[e(t)] = \gamma_e(0) + m_e^2$

Which can be generalised:

$$E[e(t)^{2}] = \gamma_{e}(0) + m_{e}^{2}$$

$$E[y(t)^{2}] = \gamma_{y}(0) + m_{y}^{2}$$

$$E[e(t)y(t-1)] = E[(e(t) - m_{e})(y(t-1) - m_{y})] + m_{y}m_{e}$$

$$E[e(t)y(t-1)] = m_{e}m_{y}$$

As the de-biased signals are incorellated!

$$\begin{split} \gamma_y(0) &= \tfrac{1}{9}(\gamma_y(0) + m_y^2) + (\gamma_e(0) + m_e^2) + \tfrac{25}{4} + \tfrac{2}{3}(m_e m_y) - \tfrac{5}{3}m_y - 5m_e = \tfrac{9}{8} \\ \text{Same computations for } \gamma_y(1), \gamma_y(2)... \end{split}$$

- Method 2: GOOD

Define two new processes:

$$\tilde{y}(t) = y(t) - \frac{9}{2} \rightarrow m_{\tilde{y}=0}$$

$$\tilde{e}(t) = e(t) - 1 \rightarrow m_{\tilde{e}=0}$$

So
$$y(t) = \tilde{y}(t) - \frac{9}{2}$$
 and $e(t) = \tilde{e}(t) + 1$:

$$\tilde{y}(t) + \frac{9}{2} = \frac{1}{3}(\tilde{y}(t-1) + \frac{9}{2}) + (\tilde{e}+1) + 2$$

$$\boxed{\tilde{y}(t) = \frac{1}{3}\tilde{y}(t-1) + \tilde{e}(t), \tilde{e} \sim WN(0, 1)}$$

Where $\tilde{y}(t)$ is the **de-biased process**

$$\gamma_{\tilde{y}(0)} = \frac{1}{1 - \frac{1}{9}} = \frac{9}{8}$$

$$\gamma_{\tilde{y}(1)} = \frac{9}{8} \frac{1}{3} = \frac{3}{8}$$

$$\gamma_{\tilde{y}(2)} = \frac{3}{8} \frac{1}{3} = \frac{1}{8} \dots$$

$$\gamma_{\tilde{y}(1)} = \frac{91}{83} = \frac{3}{8}$$

Now that we found $\gamma_{\tilde{y}(\tau)}$ we want to find $\gamma_y(\tau)$ $\gamma_y(\tau)$:

$$E[(y(t) - \frac{9}{2})(y(t-\tau) - \frac{9}{2})] = E[\tilde{y}(t)\tilde{y}(t-\tau)] = \gamma_{\tilde{y}}(\tau)$$

since $m_{\tilde{y}} = 0$ Which can be generalised :

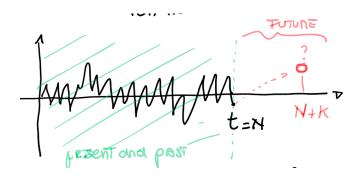
If y(t) and y(t) are two SSPs that differ only from a constant value $y(t) = \tilde{y}(t) + k$ then :

$$\boxed{\gamma_y(\tau) = \gamma_{\tilde{y}}, \forall \tau}$$

$$\boxed{\Gamma_y(w) = \Gamma_{\tilde{y}}, \forall w}$$

3 Chapter 3: Prediction

The prediction problem is to find the **best possible value** for $\hat{y}(t+k|t)$ given the **measured data** up to time t $\{y(1),...,y(N)\}$



To obtain the **optimal** prediction :

- 1. We have to make a mathematical model for $\{y(1),...,y(N)\}$
- 2. Using the model compute the optimal solution

To find the **best** mathematical model:

- 1. We select a class of models for time-series $y(t) = W(z, \theta)e(t)$ where e(t) is a WN and θ a parameter vector.
- 2. We compute the prediction of y(t) using the mathematical model:

$$\hat{y}(t+1|t;\theta)$$

3.
$$\hat{\theta} = argmin_{\theta} \left[\frac{1}{N} \sum_{t=1}^{N} (y(t) - \hat{y}(t|t-1); \theta))^2 \right]$$

4. Find $y(t)=W(Z,\hat{\theta})e(t)$ which is the best model from prediction performance . Use this to compute $\hat{y}(N+K|N)$

To create a **predictor** from an ARMA/ARMAX we need to define 2 tools:

- All-pass filter
- Canonical representation

3.1 All-Pass Filter

An All-Pass Filter is a first-order, linear, digital filter with a special constrained structure:

$$T(Z) = \frac{1}{a} \frac{Z+a}{Z+\frac{1}{a}}, a \in \Re$$

that depends on only one parameter and has a **pole** in $z = -\frac{1}{a}$ and zero in z = -a Properties :

• Magnitude

$$|T(e^{jw})|^2 = |\frac{1}{a}\frac{e^{jw}+a}{e^{jw}+\frac{1}{a}}|^2 = \frac{1}{a}(\frac{e^{jw}+a}{e^{jw}+\frac{1}{a}}) \cdot \frac{1}{a}(\frac{e^{-jw}+a}{e^{-jw}+\frac{1}{a}}) = \frac{1}{a^2}\frac{1+a^2+2acosw}{1+\frac{1}{a^2}+\frac{2cosw}{a}} = 1$$
 An all-pass filter is characterized by a **frequency response** having **unitary magnitude**:

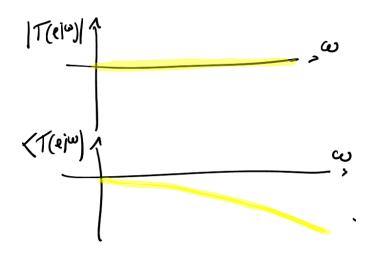
$$\Gamma_y(w) = |T(e^{jw})|^2 T_v(w)$$

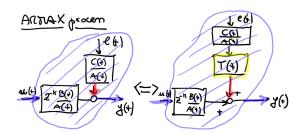
An all-pass filter does not alter the **spectrum** of its input v(t). This does **not** mean that y(t) = v(t) but that they're **statistically equivalent**:

$$\Gamma_y(w) = \Gamma_v(w)$$

$$\gamma_y(w) = \gamma_v(w)$$

Input and output are **not** identical because although there is no change in amplitude ,an all-pass filter makes a distortion in phase.





The two representations of the ARMAX process are **equivalent**. The phase distortion added to signal e(t) is **not relevant**. On the other hand, adding a T(Z) all-pass filter between u(t) and y(t) alters the behaviour of the system **critically!!**

3.2 Canonical Representation

An ARMA process can have ∞ equivalent representations (there is no way to represent it in a unique way).

There is a special representation called **Canonical Representation**:

Given a SSP y(t) that can be modellded as an ARMA process:

$$y(t) = \frac{C(Z)}{A(Z)}e(t)$$

 $\frac{C(Z)}{A(Z)}$

is the **canonical representation** if:

- 1. C(Z) and A(Z) have same degree (relative degree is 0)
- 2. C(Z) and A(Z) are **coprime** (no common factors)
- 3. C(Z) and A(Z) are **monic** (coefficient of max degree of both C(Z) and A(Z) is 1)
- 4. All roots of C(Z) and A(Z) are **strictly inside** the unit circle

Example

$$y(t) = \frac{1+3Z^{-1}}{2+Z^{-1}}e(t-2), e(t) \sim WN(0,1)$$

• Type and order

ARMA type process of order 1,3 = ARMA(1,3)

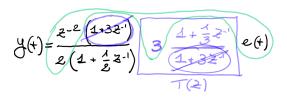
• Canonical form

$$y(t) = \frac{Z^{-2} + 3Z^{-3}}{2 + Z^{-1}}e(t)$$

- 1. Degree of C(Z) is 2, degree of A(Z) is $0 \to X$
- 2. No common factors \rightarrow OK
- 3. C(Z) is monic, A(Z) is not monic $\to X$
- 4. Zero in $-3 \rightarrow X$

By collecting and using an All-Pass Filter:

$$y(t) = \frac{Z^{-2}(1+3Z^{-1})}{2(1+\frac{1}{2}Z^{-1})} \cdot 3\frac{1+\frac{1}{3}Z^{-1}}{1+3Z^{-1}}e(t)$$



Defining $\theta = \frac{3}{2}e(t-2), \theta \sim ?$

The variance of $\theta = E[\theta(t)^2] = E[(\frac{3}{2}e(t-2))^2] = \frac{9}{4} \cdot 1$

 $\theta \sim WN(0, \frac{9}{4})? \rightarrow \Gamma_{\theta}(w) = |\frac{3}{2}e^{-2jw}|^2 \cdot 1 = \frac{9}{4}$ which is constant value so $\theta \sim WN(0, \frac{9}{4})$ Finally we obtain:

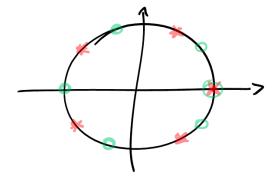
$$y(t) = \frac{1 + \frac{1}{3}Z^{-1}}{1 + \frac{1}{2}Z^{-1}}\theta(t), \theta \sim WN(0, \frac{9}{4})$$

Which is an $ARMA(1,1) \rightarrow$ the canonical representation is the representation with minimum order!

Remark

Does the canonical representation always exist?

Not always: the problem lies in the 4^{th} condition. It is possible that the system has poles or zeros **on** the unitary circle.



If C(Z) has zeros on the u.c \rightarrow prediction from data is not asymptotically stable

If A(Z) has roots in $+1 \to ARIMA$ models:

$$y(t) = \frac{C(Z)}{(Z-1)^d A(Z)} e(t) \to ARIMA(m, d, n)$$

[Autoregressive Integrated Moving average]

A special case of ARIMA \rightarrow **ARIMA**(0,1,0):

$$y(t) = \frac{1}{1 - Z^{-1}}e(t)$$

$$y(t) = y(t-1) + e(t), e(t) \sim WN(0, \lambda^2)$$

Process y(t) is a **Random Walk** that uses as TF $\frac{1}{1-Z^{-1}}$ which is an **integrator in discrete time**

ARIMA processes have an **asymptotically stable** predictor that can be used to model **not strictly stationary processes!!**

3.3 Predictor

The predictor at time t+k, given the data up to time t is:

$$\hat{y}(t+k|t)$$

The **prediction error** is:

$$\epsilon(t+k) = y(t+k) - \hat{y}(t+k|t)$$

So the **real value** is predictor + error:

$$y(t+k|t) = \hat{y}(t+k|t) + \epsilon(t+k)$$

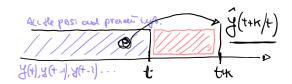
$$y(t) = \hat{y}(t|t - k) + \epsilon(t)$$

The formulas as equivalent because as per hypothesis y(t) is **stationary**.

3.3.1 Optimality

The predictor $\hat{y}(t+k|t)$ is **optimal** if:

- 1. $E[\hat{y}(t+k|t) \cdot \epsilon(t+k)] = 0$, predictor and error must be **incorrelated**
- 2. $E[y(t) \cdot \epsilon(t+k)] = E[y(t-1) \cdot \epsilon(t+k)]... = 0$



The red part shows the **unpredictable** part of y(t+k), which is the error $\epsilon(t+k)$. If error and predictor were **correlated** than some useful unused information about $\hat{y}(t+k|t)$ would be in $\epsilon(t+k)$ which means that the predictor is not optimal. The same goes for y(t), y(t-1)... of point 2): the error cannot contain information about the past/present information.

3.3.2 1-step ahead prediction of MA(n)

$$y(t) = e(t) + c_1 e(t-1) + \dots + c_n e(t-n), e(t) \sim WN(0, \lambda^2)$$

We assume the MA(n) is represented in the **canonical representation**: we must make assumptions about the 4^{th} property.

Given:

- **Present time**: $t-1 \to c_1 e(t-1) + ... + c_n e(t-n)$
- Future : $t \to e(t)$

Predictor from noise

The **optimal predictor** from **noise** is :

$$\hat{y}(t|t-1) = c_1 e(t-1) + \dots + c_n e(t-n)$$

with error

$$\epsilon(t) = y(t) - \hat{y}(t|t-1) = e(t)$$

Optimality:

•
$$E[\hat{y}(t|t-1)\epsilon(t)] = E[(c_1e(t-1) + ... + c_ne(t-n))(e(t))] = 0$$

•
$$E[y(t-1)\epsilon(t)] = E[(e(t-1)+c_1e(t-2)+...+c_ne(t-n-1))(e(t))] = ... = 0$$

Verified because of incorrelation of white noise.

Since WN is **unknown** and cannot be measured, a better predictor has to be chosen from **measurable data**.

Predictor from data

TF:

$$y(t) = (1 + c_1 Z^{-1} + \dots + c_n Z^{-n})e(t)$$

Inverse TF (Whitening Filter):

$$e(t) = \frac{1}{1 + c_1 Z^{-1} + \dots + c_n Z^{-n}} y(t)$$

$$\hat{y}(t|t-1) = (1 + c_1 Z^{-1} + \dots + c_n Z^{-n})e(t) \to \hat{y}(t|t-1) = \frac{c_1 Z^{-1} + \dots + c_n Z^{-n}}{1 + c_1 Z^{-1} + \dots + c_n Z^{-n}}y(t)$$

Collecting Z^{-1} :

$$\hat{y}(t|t-1) = \frac{c_1 + \dots + c_n Z^{-n+1}}{1 + c_1 Z^{-1} + \dots + c_n Z^{-n}} y(t-1)$$

As seen in time-domain representation the prediction makes use of **present** and past data as well as past predictions.

3.3.3 K-steps ahead predictor of MA(n)

$$y(t) = e(t) + c_1 e(t-1) + \dots + c_{k-1} e(t-k+1) + c_k e(t-k) + \dots + c_n e(t-n)$$

Given:

-Present time: $k \rightarrow c_k e(t-k) + ... + c_n e(t-n)$

-Future: $t \to e(t) + ... + c_{k-1}e(t-k+1)$

Predictor from noise

$$\hat{y}(t|t-k) = c_k e(t-k) + ... + c_n e(t-n)$$

with error:

$$\epsilon(t) = e(t) + \dots + c_{k-1}e(t-k+1)$$

Predictor from data

$$\hat{y}(t|t-k) = \frac{c_k + c_{k+1}Z^{-1} + \dots + c_nZ^{-n+k}}{1 + c_1Z^{-1} + \dots + c_nZ^{-n}} y(t-k)$$

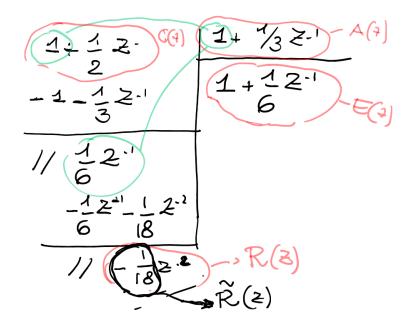
3.3.4 K-steps ahead predictor of general ARMA(m,n)

$$y(t) = \frac{C(Z)}{A(Z)}e(t), e(t) \sim WN(0, \lambda^2)$$

(Under the hypothesis of canonical representation)

The AR(m) part presents a recursion : need to introduce k-steps **polynomial** division between C(Z) and A(Z) obtaining :

- $\mathbf{E}(\mathbf{Z}) \to \text{result (quotient)}$
- $\mathbf{R}(\mathbf{Z}) \to \text{residual (remainder)}$



$$C(Z) = E(Z)A(Z) + R(Z)$$
$$\frac{C(Z)}{A(Z)} = E(Z) + \frac{R(Z)}{A(Z)}$$

Noting that in k-steps division R(Z) can be rewritten by collecting Z^{-k} :

$$R(Z) = Z^{-k}\tilde{R}(Z)$$

$$\boxed{\frac{C(Z)}{A(Z)} = E(Z) + \frac{Z^{-k}\tilde{R}(Z)}{A(Z)}}$$

The new transfer function is:

$$y(t) = [E(Z) + \frac{Z^{-k}\tilde{R}(Z))}{A(Z)}]e(t)$$

$$y(t) = E(Z)e(t) + \frac{\tilde{R}(Z)}{A(Z)}e(t-k)$$

Where E(Z)e(t) is the **unpredictable part** as it depends on e(t),...,e(t-k+1)

Predictor from noise

$$\hat{y}(t|t-k) = \frac{\tilde{R}(Z)}{A(Z)}e(t-k)$$

with error:

$$\epsilon(t) = E(Z)e(t)$$

Predictor from data

$$y(t) = \frac{C(Z)}{A(Z)}e(t) \xrightarrow{\text{Whitening}} e(t) = \frac{A(Z)}{C(Z)}y(t)$$
$$\hat{y}(t|t+k) = \frac{\tilde{R}(Z)Z^{-k}}{A(Z)} \cdot \frac{A(Z)}{C(Z)}y(t)$$
$$\hat{y}(t|t-k) = \frac{\tilde{R}(Z)}{C(Z)}y(t-k)$$

Remark 1

Both the predictor from noise and data work under the assumption of SSP. The stationary property is satisfied if both A(Z) and C(Z) have all roots (poles) strictly inside the unitary circle. But this is satisfied by the 4^{th} condition of the canonical representation hypothesis.

Remark 2

$$\epsilon(t) = y(t) - \hat{y}(t|t-k) = E(Z)e(t)$$
 where $E(Z)$ is a SSP of type $\mathbf{MA(k-1)}$

Remark 3

In the case of K=1 the polynomial division result in :

 $-\mathbf{E}(\mathbf{Z}) = 1$ as both $C(\mathbf{Z}), A(\mathbf{Z})$ are monic and have same degree

$$-\mathbf{R}(\mathbf{Z}) = \mathbf{C}(\mathbf{Z}) - \mathbf{A}(\mathbf{Z})$$

which results in

$$\hat{y}(t|t-k) = \frac{C(Z) - A(Z)}{A(Z)}e(t)$$

$$\epsilon(t) = e(t)$$

Instead of having term R(Z) , the formula is now C(Z)-A(Z). As $R(Z) = \tilde{R}(Z)Z^{-1}$ there is a hidden Z^{-1} in C(Z)-A(Z).

3.3.5 K-steps ahead prediction of ARMAX(m,n,k+p)

$$y(t) = \frac{B(Z)}{A(Z)}u(t-k) + \frac{C(Z)}{A(Z)}e(t), e(t) \sim WN(0, \lambda^2)$$

Where:

$$A(Z) = 1 + a_1 Z^{-1} + \dots + a_m Z^{-m}$$

$$B(Z) = b_0 + b_1 Z^{-1} + \dots + b_p Z^{-p}$$

$$C(Z) = 1 + c_1 Z^{-1} + \dots + c_n Z^{-n}$$

In the hypothesis that $\frac{C(Z)}{A(Z)}$ is in **canonical representation** and keeping in mind that for $\frac{B(Z)}{A(Z)}u(t-k)$ no **spectral equivalence** modifications can be made. In an ARMAX(m,n,k+p) process the most interesting prediction that can be made is the **delay** between u(t) and y(t) \rightarrow **k** so we'll deal only with k-steps predictions.

Predictor from noise

Separate predictable from unpredictable part in $\frac{C(Z)}{A(Z)}e(t)$ K-steps division $\frac{C(Z)}{A(Z)} \to E(Z) + \frac{R(Z)}{A(Z)}$

$$y(t) = \frac{B(Z)}{A(Z)}u(t-k) + E(Z)e(t) + \frac{\tilde{R}(Z)}{A(Z)}e(t-k)$$

where

$$\frac{B(Z)}{A(Z)}u(t-k) \to \text{depends on } u(t-k),...,u(t-k-p) \to \text{predictable}$$

 $E(Z)e(t) \rightarrow \text{depends on } e(t), e(t-1), ..., e(t-k+1) \rightarrow \text{unpredictable}$

$$\frac{\tilde{R}(Z)}{A(Z)}e(t-k) \to \text{depends on } e(t-k),...,e(t-k-p) \to \text{predictable}$$

so

$$\hat{y}(t|t-k) = \frac{B(Z)}{A(Z)}u(t-k) + \frac{R(Z)}{A(Z)}e(t)$$

$$\epsilon(t) = y(t) - \hat{y}(t|t - k) = E(t)e(t)$$

Which is optimal if

- $\epsilon(t) \perp \hat{y}(t|t-k)$
- $\epsilon(t) \perp y(t-k), y(t-k-1)...$

Predictor from data

$$e(t) = \frac{A(Z)}{C(Z)}y(t) - \frac{B(Z)}{A(Z)}u(t-k)$$

$$\hat{y}(t|t-k) = \frac{B(Z)}{A(Z)}u(t-k) + \frac{R(Z)}{A(Z)}\left[\frac{A(Z)}{C(Z)}y(t) - \frac{B(Z)}{A(Z)}u(t-k)\right]$$

$$\hat{y}(t|t-k) = \frac{R(Z)}{C(Z)}y(t) + \left[\frac{B(Z)}{A(Z)} - \frac{R(Z)B(Z)}{A(Z)C(Z)}\right]u(t-k)$$

$$\hat{y}(t|t-k) = \frac{B(Z)}{C(Z)}y(t) + \left[\frac{B(Z)(C(Z)-R(Z))}{A(Z)C(Z)}\right]u(t-k)$$
Knowing that $C(Z) = A(Z)E(Z) + R(Z) \rightarrow C(Z) - R(Z) = A(Z)E(Z)$

$$\hat{y}(t|t-k) = \frac{B(Z)E(Z)}{C(Z)}u(t-k) + \frac{\tilde{R}(Z)}{C(Z)}y(t-k)$$

$$\hat{x}(t|t-k) = \frac{B(Z)E(Z)}{C(Z)}u(t-k) + \frac{\tilde{R}(Z)}{C(Z)}y(t-k)$$

Note that $\frac{\tilde{R}(Z)}{C(Z)}y(t-k)$ is the exact ARMA predictor.

The prediction error is the same as in the **ARMA** process: the **exogenous** part does not add any **additional uncertainty**.

Remark: Special case k=1

$$y(t) = \frac{B(Z)}{A(Z)}u(t-1) + \frac{C(Z)}{A(Z)}e(t)$$

$$E(Z) = 1 \text{ and } R(Z) = C(Z) - A(Z)$$

$$\hat{y}(t|t-1) = \frac{B(Z)}{C(Z)}u(t-1) + \frac{C(Z) - A(Z)}{C(Z)}y(t)$$

3.4 Examples & Exercises

3.4.1 Example 1

Given a process

$$y(t) = \frac{Z+3}{2Z+1}e(t-1), e(t) \sim WN(0,1)$$

Since the pole of the TF is $z = -\frac{1}{2}$ inside the unitary circle, W(Z) is asymptotically stable \rightarrow y(t) is **stationary**.

1. Compute $\gamma_y(0)$

NB.: To calculate the variance it is not important for the system to be in canonical representation

 $y(t) = \frac{C(Z)}{A(Z)}e(t-1)$ is **not canonical** since it has

- Z = -3 not inside unitary circle
- 2Z in the A(Z) term
- $-Z^{-1}e(t)$

Using an All-Pass Filter:

$$y(t) = \frac{Z+3}{2(Z+\frac{1}{2})}Z^{-1} \cdot 3\frac{Z+\frac{1}{3}}{Z+3}e(t)$$
$$\eta = \frac{3}{2}Z^{-1}e(t) \sim WN(0, \frac{9}{4})$$
$$y(t) = \frac{Z+\frac{1}{3}}{Z+\frac{1}{2}}\eta(t)$$

Passing in time domain:

$$y(t) = -\frac{1}{2}y(t-1) + \eta(t) + \frac{1}{3}\eta(t-1)$$

$$-m_y = E[y(t)] = -\frac{1}{2}E[y(t-1)] + \frac{4}{3}m_e \to 0$$

$$-\gamma_y(0) = E[y(t)^2] = E[(-\frac{1}{2}y(t-1) + \eta(t) + \frac{1}{3}\eta(t-1))^2]$$

$$\gamma_y(0) = \frac{1}{4}\gamma_y(0) + \frac{9}{4} + \frac{1}{9}\frac{9}{4} - \frac{1}{3}E[y(t-1)\eta(t-1)]$$

$$\frac{3}{4}\gamma_y(0) = \frac{10}{4} - \frac{1}{3}E[(-\frac{1}{2}y(t-2) + \eta(t-1) + \frac{1}{3}\eta(t-2))\eta(t-1)]$$

$$\frac{3}{4}\gamma_y(0) = \frac{10}{4} - \frac{1}{3}E[\eta(t-1)^2] \to \frac{3}{4}\gamma_y(0) = \frac{10}{4} - \frac{1}{3}\frac{9}{4}$$

$$\boxed{\gamma_y(0) = \frac{7}{3}}$$

2. Prediction for k=1

Using the canonical negative power representation

$$y(t) = \frac{1 + \frac{1}{3}Z^{-1}}{1 + \frac{1}{2}Z^{-1}}\eta(t)$$

Applying the k=1 prediction formula $\hat{y}(t|t-1) = \frac{C(Z) - A(Z)}{C(Z)}y(t)$:

$$\hat{y}(t|t-1) = \frac{1 + \frac{1}{3}Z^{-1} - 1 - \frac{1}{2}Z^{-1}}{1 + \frac{1}{3}Z^{-1}}y(t)$$

$$\hat{y}(t|t-1) = \frac{-\frac{1}{6}}{1 + \frac{1}{3}Z^{-1}}y(t-1)$$

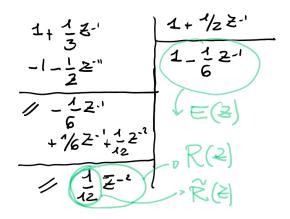
In time domain:

$$\hat{y}(t|t-1) = -\frac{1}{3}\hat{y}(t-1|t-2) - \frac{1}{6}y(t-1)$$

$$\epsilon(t) = y(t) - \hat{y}(t|t-1) = E(Z)\eta(t) = \eta(t)$$

$$\boxed{var[y(t) - \hat{y}(t|t-1)] = var[\eta(t)] = \frac{9}{4}}$$

3. Prediction for k=2



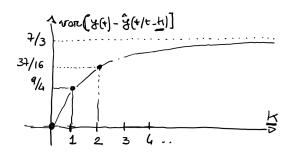
$$\hat{y}(t|t-2) = \frac{R(Z)}{C(Z)}y(t) = \frac{\tilde{R}(Z)}{C(Z)}y(t-2)$$
$$\hat{y}(t|t-2) = \frac{\frac{1}{12}}{1 + \frac{1}{3}Z^{-1}}y(t-2)$$

$$\hat{y}(t|t-2) = -\frac{1}{3}\hat{y}(t-1|t-3) + \frac{1}{12}y(t-2)$$

$$\epsilon(t) = y(t) - \hat{y}(t|t-2) = E(Z)\eta(t) = (1 - \frac{1}{6}Z^{-1})\eta(t)$$

$$var[y(t) - \hat{y}(t|t-2)] = var[(1 - \frac{1}{6}Z^{-1})\eta(t)] = \frac{37}{16}$$

4. Properties of $var[\epsilon(t)]$ as function of k



$$-k = 0 \rightarrow var[y(t) - \hat{y}(t|t-k)] = 0$$

-
$$k = 1 \rightarrow var[y(t) - \hat{y}(t|t-k)] = \lambda^2$$

- $k \to \infty \to var[y(t) - \hat{y}(t|t-k)] = \gamma_y(0)$ because when $k \to \infty$ the prediction goes to zero!

- $var[y(t) - \hat{y}(t|t-k)]$ is a monotonic (not strictly) increasing function

5. Prediction goodness

The Error to signal ratio is a useful prediction measure:

$$ESR(k) = \frac{var[y(t) - \hat{y}(t|t - k)]}{var[y(t)]}$$

For k=1

$$ESR(1) = \frac{\frac{9}{4}}{\frac{7}{3}} = 0.97$$

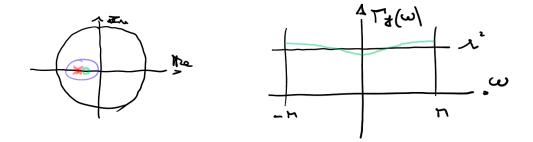
Which is a very bad prediction. The most trivial prediction that can be done is

$$\hat{y}(t|t-k) = m_y$$

(predicting the mean) which has ESR(k)=1.

For k=1 we only have a 3% better prediction than the trivial one.

The predictor for k=1 is **optimal** which means that **no better** prediction can be made: the bad prediction is an intrinsic property of the process y(t).



By analysing the poles and zeros, it is easy to see that they're so close together that they almost cancel each other out.

The **spectrum** $\Gamma_y(w)$ in green is very close to that of the **white noise**: this is the reason y(t) is hard to predict

3.4.2 Example 2 - Practical

We have measured 5 data points of a signal:

$$y(1) = 1, y(2) = \frac{1}{2}, y(3) = -\frac{1}{2}, y(4) = 0, y(5) = -\frac{1}{2}$$

With t=5 represent the present time, make a prediction of $\hat{y}(6|5)$. To solve the problem we must make a mathematical modelling assumption. Since we're still not able to do this we need some interpretations models for this signal

Model A

$$y(t) = \frac{1}{2}y(t-1) + \frac{1}{4}y(t-2) + e(t), e(t) \sim WN(0, \lambda^2)$$

Model B

$$y(t) = e(t) + \frac{1}{2}e(t-1), e(t) \sim WN(0, \lambda^2)$$

To determine which model is better we compute the **optimal** model assuming the chosen model is right.

• Assuming Model A right

$$y(t) = \frac{1}{1 - \frac{1}{2}Z^{-1} - \frac{1}{4}Z^{-2}}e(t)$$

is an AR(2) process in canonical representation (check it always!).

Since we're dealing with a k=1 prediction :

$$\hat{y}(t|t-1) = \frac{C(Z) - A(Z)}{C(Z)}y(t)$$

$$\hat{y}(t|t-1) = \frac{1 - 1 + \frac{1}{2}Z^{-1} + \frac{1}{4}Z^{-2}}{1}y(t)$$
$$\hat{y}(t|t-1) = \frac{1}{2}y(t-1) + \frac{1}{4}y(t-2)$$

Substituting the data points:

$$\hat{y}(6|5) = \frac{1}{2}y(5) + \frac{1}{4}y(4) = -\frac{1}{4}$$

• Assuming Model B right

$$y(t) = (1 + \frac{1}{2}Z^{-1})e(t)$$

is an MA(1) process in canonical representation.

Since we're dealing with a k=1 prediction :

$$\hat{y}(t|t-1) = \frac{C(Z) - A(Z)}{C(Z)} y(t)$$

$$\hat{y}(t|t-1) = \frac{1 + \frac{1}{2}Z^{-1} - 1}{1 + \frac{1}{2}Z^{-1}} y(t)$$

$$\hat{y}(t|t-1) = -\frac{1}{2}\hat{y}(t-1|t-2) + \frac{1}{2}y(t-1)$$

Substituting the data points:

$$\hat{y}(6|5) = -\frac{1}{2}\hat{y}(5|4) + \frac{1}{2}y(5) = -\frac{1}{2}\hat{y}(5|4) - \frac{1}{4}$$

To compute $-\frac{1}{2}\hat{y}(5|4)$ we need to go back to the **initial condition** to compute all terms up to time =5:

$$-\hat{y}(2|1)=-\frac{1}{2}\hat{y}(1|0)+\frac{1}{2}y(1)$$
 by making the assumption that $-\frac{1}{2}\hat{y}(1|0)=m_y\to\frac{1}{2}$

$$-\hat{y}(3|2) = -\frac{1}{2}\hat{y}(2|1) + \frac{1}{2}y(2) = 0$$

$$-\hat{y}(4|3) = -\frac{1}{2}\hat{y}(3|2) + \frac{1}{2}y(3) = -\frac{1}{4}$$

$$-\hat{y}(5|4) = -\frac{1}{2}\hat{y}(4|3) + \frac{1}{2}y(4) = \frac{1}{8}$$

$$-\hat{y}(6|5) = -\frac{1}{2}\hat{y}(5|4) + \frac{1}{2}y(5) = -\frac{5}{16}$$

Our final prediction for model B is $\hat{y}(6|5) = -\frac{5}{16}$ which depends on the initial condition made assuming that $-\frac{1}{2}\hat{y}(1|0) = m_y$. It the choice of the initial condition important? If the system is asymptotically stable, and N is big the initial condition is not important as it will vanish.

3.4.3 Example 3 - ARMAX & ARX

$$y(t) = (Z + 6Z^{-1})u(t - 2) + \frac{2}{3 + \frac{3}{2}Z - 1}\eta(t - 1), \eta \sim WN(0, 1)$$

Find predictor from data and the corresponding error with its variance. $u(t-2) \rightarrow k=2$

Canonical form for ARMA part

$$\frac{2}{3(1+\frac{1}{2}Z^{-1})}Z^{-1}\eta(t)$$

$$=\frac{2}{2}\eta(t-1) e(t) \approx WN(t)$$

So

$$e(t) = \frac{2}{3}\eta(t-1), e(t) \sim WN(0, \frac{4}{9})$$
$$\frac{1}{1 + \frac{1}{2}Z^{-1}}e(t)$$

Substituting in the original process:

$$y(t) = (Z + 6Z^{-1})u(t - 2) + \frac{1}{1 + \frac{1}{2}Z^{-1}}e(t), e(t) \sim WN(0, \frac{4}{9})$$

Is the term $(Z + 6Z^{-1})u(t - 2)$ also in canonical representation? Wrong question, there is nothing we can do about it!

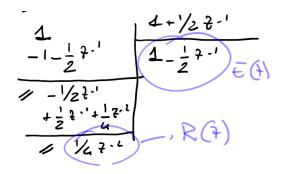
We need the form:

$$y(t) = \frac{B(Z)}{A(Z)}u(t-k) + \frac{C(Z)}{A(Z)}e(t)$$

So rewriting:

$$y(t) = \frac{(2+6Z-1)(1+\frac{1}{2}Z^{-1})}{(1+\frac{1}{2}Z^{-1})}u(t-2) + \frac{1}{1+\frac{1}{2}}e(t)$$

Using a k-steps long division $\frac{C(Z)}{A(Z)}$:



$$\hat{y}(t|t-2) = \frac{(2+6Z^{-1})(1+\frac{1}{2}Z^{-1})(1-\frac{1}{2}Z-1)}{1}u(t-2) + \frac{\frac{1}{4}Z^{-2}}{1}y(t)$$

$$\hat{y}(t|t-2) = 2u(t-2) + 6u(t-3) - \frac{1}{2}u(t-4) - \frac{3}{2}u(t-5) + \frac{1}{4}y(t-2)$$

No old prediction is used \rightarrow process is ARMAX(1,0,2+2) \rightarrow ARX(1,4) model.

$$\epsilon = E(Z)e(t) = (1 - \frac{1}{2})Z^{-1}e(t)$$

The variance of ϵ :

$$var[\epsilon(t)] = (1 + \frac{1}{4}) \cdot \frac{4}{9} = \frac{5}{4} \cdot \frac{4}{9} = \frac{5}{9}$$

3.4.4 Example ARMA with non-zero mean

$$y(t) = e(t) + 4e(t-1), e \sim WN(1,1)$$

Compute $\hat{y}(t|t-1)$ and $\hat{y}(t|t-2)$ from data.

Canonical form representation

$$y(t) = (1 + 4Z^{-1})e(t) \to y(t) = (1 + 4Z^{-1})\left[4 \cdot \frac{1 + \frac{1}{4}Z^{-1}}{1 + 4Z^{-1}}\right]e(t)$$

Getting the new $\eta(t)$:

$$\eta(t) = 4e(t)$$

$$m_{\eta} = E[\eta(t)] = E[4e(t)] = 4$$

$$var[\eta] = E[(\eta(t) - 4)^2] = E[(4e(t) - 4)^2] = 16E[(e(t) - 1)^2] = 16$$

Canonical form:

$$y(t) = (1 + \frac{1}{4}Z^{-1})\eta(t)$$

$$\eta(t) \sim WN(4, 16)$$

Method 1

De-biasing technique:

$$\tilde{y}(t) = y(t) - m_y$$

$$\tilde{\eta}(t) = \eta(t) - m_n$$

Mean of y:

$$E[y(t)] = E[(\eta(t) + \frac{1}{4}\eta(t-1))] \to m_y = \frac{5}{4}m_\eta = 5$$

So:

$$\tilde{y}(t) = y(t) - 5$$

$$\tilde{\eta}(t) = \eta(t) - 4 \to \tilde{\eta} \sim WN(0, 16)$$

Obtaining:

$$\tilde{y} + 5 = (\tilde{\eta}(t) + 4) + \frac{1}{4}(\tilde{\eta}(t-1) + 4)$$
$$\tilde{y} = \tilde{\eta}(t) + \frac{1}{4}\tilde{\eta}(t-1)$$

Now we can compute the predictions for \tilde{y} for k=1:

$$\hat{\tilde{y}}(t|t-1) = \frac{(1 + \frac{1}{4}Z^{-1}) - 1}{(1 + \frac{1}{4}Z^{-1})}\tilde{y}(t)$$

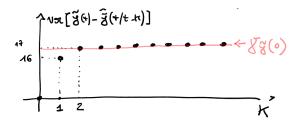
$$\hat{\tilde{y}}(t|t-1) = \frac{\frac{1}{4}}{(1 + \frac{1}{4}Z - 1)}\tilde{y}(t-1)$$
$$\epsilon(t) = \tilde{\eta}(t) = 16$$

Now we can compute the prediction for \tilde{y} for k=2:

Which means that

$$\hat{\tilde{y}}(t|t-2) = 0$$

Because MA(1) process has a **finite memory** of 1-step only! So $\tilde{\epsilon}(t) = \tilde{y}(t) - \hat{\tilde{y}}(t|t-2) = \tilde{y}(t)$ so the $var[\tilde{\epsilon}(t)] = var[\tilde{y}(t)] = (1 + \frac{1}{16}) \cdot 16 = 17$



We need to go back to the original process because $\hat{y}(t|t-1) \neq \hat{\hat{y}}(t|t-1)$: - for **k=1**

$$\hat{y}(t|t-1) - 5 = \frac{\frac{1}{4}}{1 + \frac{1}{4}Z^{-1}}(y(t-1) - 5)$$

$$\hat{y}(t|t-1) = \frac{\frac{1}{4}}{1 + \frac{1}{4}Z^{-1}}y(t-1) + 5 - 5\frac{\frac{1}{4}}{1 + \frac{1}{4}Z^{-1}}$$

The term $5\frac{\frac{1}{4}}{1+\frac{1}{4}Z^{-1}}$ can be resolved by applying the **frequency response** theorem by taking in account that 5 is a sinusoid with frequency w=0:

$$\frac{\frac{1}{4}}{1 + \frac{1}{4}e^{0j}} \cdot 5 = 1$$

so:

$$\hat{y}(t|t-1) = \frac{\frac{1}{4}}{1 + \frac{1}{4}Z^{-1}}y(t-1) + 4$$

- for k=2

$$\hat{\tilde{y}}(t|t-2) = 0 \rightarrow \hat{y}(t|t-2) - 5 = 0$$

$$\hat{y}(t|t-2) = 5$$

$$var[\tilde{\eta}(t)] = var[\eta(t)]$$

Method 2

De-bias technique only on $\eta(t)$:

$$\tilde{\eta}(t) = \eta(t) - 4$$
$$y(t) = \tilde{\eta}(t) + 4 + \frac{1}{4}(\tilde{\eta}(t) + 4)$$

$$y(t) = \tilde{\eta}(t) + \frac{1}{4}\tilde{\eta}(t-1) + 5$$

Which can be considered as an **ARMAX** process.

-for k=1

$$y(t) = u(t-1) + (1 + \frac{1}{4}Z^{-1}\tilde{\eta}(t)), u(t) = 5 \forall t$$

u(t-1) is chosen arbitrarily because we need to compute $\hat{y}(t|t-1)$:

- k = 1
- -B(Z) = 1
- $-C(Z) = 1 + \frac{1}{4}Z^{-1}$
- -A(Z) = 1

$$\hat{y}(t|t-1) = \frac{1 \cdot 1}{1 + \frac{1}{4}Z^{-1}}u(t-1) + \frac{\left(1 + \frac{1}{4}Z^{-1}\right) - 1}{1 + \frac{1}{4}Z^{-1}}y(t)$$

As u(t-1) = 5 the system taking it as input can be simplified again using FR. Theorem :

$$\hat{y}(t|t-1) = \frac{\frac{1}{4}}{1 + \frac{1}{4}Z^{-1}}y(t-1) + 4$$

Which is the same result as for the first method.

-for k=2

$$y(t) = u(t-2) + (1 + \frac{1}{4}Z^{-1})\tilde{\eta}(t), u(t) = 5\forall t$$

Making a 2 step long division $\frac{C(Z)}{A(Z)}$:

$$-E(Z) = 1 + \frac{1}{4}Z^{-1} - R(Z) = 0$$

$$\hat{y}(t|t-2) = \frac{1 \cdot \left(1 + \frac{1}{4}Z^{-1}\right)}{1 + \frac{1}{4}Z^{-1}}u(t-2) + 0$$

$$\hat{y}(t|t-2) = 5$$

Which again is the same result as in the first method.

4 Chapter 4: Identification

The focus of MIDA 1 are **parametric** identification or learning techniques. They are the most used and popular identification techniques but many non-parametric techniques are essential for identification (ex: state-space identification, spectrum estimation, unsupervised learning...)

Any parametric identification technique is based on a five step approach:

1. Experiment design & data collection

This step deals with **designing** the experiment, selecting the **length N** of the dataset and **data pre-processing**.

2. Selection of a class of parametric models

This steps deals with the selection of **class** of parametric models $m(\theta)$ where θ is the unknown parameter vector. Our focus will be on :

- -discrete time
- -dynamic
- -linear
- -time-invariant

systems. As already seen **ARMAX & ARMA** are the most general class of models for these systems.

3. Selection of a performance index

A function $J(\theta) \geq 0$ that tells the **ordering** of different models. The performance index assesses the **quality** of a model.

The **prediction error method** is the choice for our performance index:

$$J_N(\theta) = \frac{1}{N} \sum_{t=1}^{N} (y(t) - \hat{y}(t|t-1,\theta))^2$$

that represents the **sample variance** of the prediction error computed on the available dataset of length N.

The P.E.M assumes that the ability of a model to make a good prediction of the future is a good quality index for the model.

4. Optimization

Optimization consists in **minimizing** $J(\theta)$ with respect to θ :

$$\hat{\theta} = argmin_{\theta} \{ J(\theta) \}$$

so that

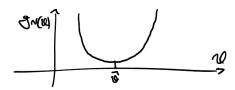
$$m(\hat{\theta})$$

is the **optimal model** on the class of models $m(\theta)$.

$$J_N(\theta) = R^{n_\theta} \to R^+$$

In optimisation 3 different situations can be found:

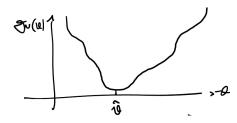
• $J(\theta)$ is **quadratic** function of θ



 J_N is a quadratic function of θ : in this case it's usually easy to find the global minimum explicitly.

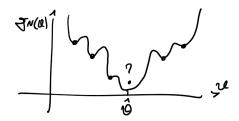
AR & ARX models are of this kind.

• $J(\theta)$ is **not** a quadratic function , **no local minima**



In this case the function has no local minima so the **unique solution** is guaranteed to be found using an **iterative algorithm**.

• $J(\theta)$ is **not** quadratic, with local minima



In this case the function has local minima so using an **iterative algorithm** is the best way to find the unique solution which is **not guaranteed** to be found.

ARMAX & ARMA models are of this kind.

5. Validation

The validation step checks if $m(\hat{\theta})$ can be considered a **valid** model for our purposes. Usually a technique called **cross-validation** is used.

4.1 Identification of ARX models

Given an available dataset of length N:

$$\{u(1), u(2), ..., u(N)\}$$

$${y(1), y(2), ..., y(N)}$$

An the model class ARX(m,p+1):

$$y(t) = \frac{B(Z)}{A(Z)}u(t-1) + \frac{1}{A(Z)}e(t), e(t) \sim WN(0, \lambda^2)$$

where $\theta = [a_1...a_m b_0...b_p]^T$ is the **parameter vector** of dimension $n_\theta = m + p + 1$.

Remark

Using k=1 is not a restriction but the **most general** case of an ARX. If the system has k > 1 we will find out during the identification process.

4.1.1 Loss function: Least Squares

The loss function for the ARX models is the **P.E.M.**:

$$J_N(\theta) = \frac{1}{N} \sum_{t=1}^{N} (y(t) - \hat{y}(t|t-1,\theta))^2$$

The predictor for the model, deriving it from the general ARMAX model, is:

$$\hat{y}(t|t-1;\theta) = \frac{B(Z)}{u}(t-1) + (1 - A(Z))y(t)$$

$$\hat{y}(t|t-1;\theta) = b_0 u(t-1) + \dots + b_p u(t-p-1) - a_1 y(t-1) - \dots - a_m y(t-m)$$

where we can define the **data vector**:

$$\phi = \left[-y(t-1), \dots - y(t-m), u(t-1), \dots u(t-p-1) \right]^{T}$$

so $\hat{y}(t|t-1) = \phi(t)^T \theta$, a linear function of θ . Substituting in the loss function:

$$J_N(\theta) = \frac{1}{N} \sum_{t=1}^{N} (y(t) - \phi(t)\theta)^2$$

A quadratic function is obtained so the unique solution can be found explicitly using a minimization method. To find the minimum we differentiate wrt to the parameter vector θ :

$$\frac{\partial J_N(\theta)}{\partial \theta} = 0$$

$$\frac{\partial J_N(\theta)}{\partial \theta} = \frac{2}{N} \sum_{t=1}^N \phi(t) (y(t) - \phi(t)^T \theta)$$

$$(\sum_{t=1}^N \phi(t) \phi(t)^T) \theta = \sum_{t=1}^N y(t) \phi(t)$$

Assuming that the $n_{\theta} \ge n_{\theta}$ matrix $\sum_{t=1}^{N} \phi(t)\phi(t)^{T}$ matrix is **non singular** and thus **invertible**:

$$\hat{\theta}_N = (\sum_{t=1}^N \phi(t)\phi(t)^T)^{-1} (\sum_{t=1}^N y(t)\phi(t))$$

This is the **explicit** solution of the ARX identification problem also known as **Least Squares**

4.1.2 Example

Consider a dataset of length N=10 and

$$y(t) = \frac{b}{1 + aZ^{-1}}u(t - 1) + \frac{1}{1 + aZ^{-1}}e(t), e(t) \sim WN(0, \lambda^2)$$

an ARX(1,1) general model class. Assuming that the process is in canonical representation (|a| < 1 must hold). The predictor of the model is:

$$\hat{y}(t|t-1) = \frac{B(Z)}{1}u(t-1) + \frac{1 - A(Z)}{1}e(t)$$

$$\hat{y}(t|t-1) = bu(t-1) - ay(t-1)$$

and $\theta = [a, b]^T$.

Method 1

The loss function is:

$$J_{10}(\theta) = \frac{1}{10} \sum_{t=1}^{10} (y(t) - bu(t-1) + ay(t-1))^2$$

Remark

Since we don't have data points for t=0 , starting at time t=1 doesn't allow us to compute -bu(0) + ay(0) so a modified version of the performance index is used:

$$J_N(\theta) = \frac{1}{N-h} \sum_{t=h+1}^{N} (y(t) - \hat{y}(t|t-1))^2$$

where $h = max\{m, p+1\}$

In our case h = max(1,1) = 1 so

$$J_{10}(\theta) = \frac{1}{9} \sum_{t=2}^{10} (y(t) - bu(t-1) + ay(t-1))^2$$

To obtain the best parameter vector $\rightarrow \frac{\partial J_{10}(\theta)}{\partial \theta}=0$ where $\theta=[a,b]^T$:

$$\frac{\partial J_{10}(\theta)}{\partial \theta} = \begin{cases} \frac{\partial J_{10}(\theta)}{\partial a} = \frac{2}{9} \sum_{t=2}^{10} (y(t) - bu(t-1) + ay(t-1)) \cdot y(t-1) = 0\\ \frac{\partial J_{10}(\theta)}{\partial b} = \frac{2}{9} \sum_{t=2}^{10} (y(t) - bu(t-1) + ay(t-1)) \cdot (-u(t-1)) = 0 \end{cases}$$

Which can be rewritten as:

$$\begin{bmatrix} \sum_{t=2}^{10} y(t-1)^2 & -\sum_{t=2}^{10} y(t-1)u(t-1) \\ -\sum_{t=2}^{10} y(t-1)u(t-1) & \sum_{t=2}^{10} u(t-1)^2 \end{bmatrix} \cdot \begin{bmatrix} a \\ b \end{bmatrix} = \begin{bmatrix} -\sum_{t=2}^{10} y(t-1)y(t) \\ \sum_{t=2}^{10} y(t)u(t-1) \end{bmatrix}$$

so

$$\begin{bmatrix} \hat{a}_{10} \\ \hat{b}_{10} \end{bmatrix} = \begin{bmatrix} \sum_{t=2}^{10} y(t-1)^2 & -\sum_{t=2}^{10} y(t-1)u(t-1) \\ -\sum_{t=2}^{10} y(t-1)u(t-1) & \sum_{t=2}^{10} u(t-1)^2 \end{bmatrix}^{-1} \begin{bmatrix} -\sum_{t=2}^{10} y(t-1)y(t) \\ \sum_{t=2}^{10} y(t)u(t-1) \end{bmatrix}$$

Method 2

Consider the predictor $\hat{y}(t|t-1) = bu(t-1) + ay(t-1)$ and the available data. Assuming that the predictor makes the **perfect prediction** on the measured data:

$$\begin{cases}
-ay(1) + bu(1) = y(2) \\
-ay(2) + bu(2) = y(3) \\
... \\
-ay(9) + bu(9) = y(10)
\end{cases}$$

Which can be separated into two matrices:

$$\Phi = \begin{bmatrix} -y(1) & u(1) \\ -y(2) & u(2) \\ \dots & \dots \\ -y(9) & u(9) \end{bmatrix} Y = \begin{bmatrix} y(2) \\ y(3) \\ \dots \\ y(10) \end{bmatrix}$$

Obtaining a **linear** system of 9 equations and 2 unknowns:

$$\Phi \cdot \theta = Y$$

Remark

• Undetermined linear system

Number of unknowns ξ number of equations \rightarrow infinite solutions

• Square linear system

Number of unknowns = number of equations \rightarrow 1! solution

• Over determined linear system

Number of unknowns; number of equations \rightarrow No solutions!

Unfortunately our case is the last one with no solutions! In this case a **least** squares (approximate) solution can be found using the square matrix : $\Phi\theta = Y \to \Phi^T \Phi\theta = \Phi^T Y$

$$\hat{\theta} = [\Phi^T \Phi]^{-1} \Phi^T Y$$

where $[\Phi^T \Phi]^{-1} \Phi^T = \Phi^+$ is the **pseudo inverse** of Φ . By creating the matrices above we obtain the **same** result for $\hat{\theta}$

4.1.3 Example 2

Assume a dataset of 5 points collected from an SSP with zero mean. $y(1) = \frac{1}{2}, y(2) = 0, y(3) = -1, y(4) = -\frac{1}{2}, y(5) = \frac{1}{4}$ and we want to make a prediction $\hat{y}(6|5)$

1. Build the model

We selected a class of model AR(1) (reason discussed further on).

$$y(t) = ay(t-1) + e(t)e(t) \sim WN(0, \lambda^2)$$

is our $m(\theta)$. Transforming $\to y(t) = \frac{1}{1-aZ^{-1}}e(t)$ which is in canonical representation if |a<1|. The corresponding k=1 predictor is

$$\hat{y}(t|t-1) = \frac{C(Z) - A(Z)}{C(Z)}y(t) \to \hat{y}(t|t-1) = ay(t-1)$$

with **performance index**:

$$J_5(a) = \frac{1}{4} \sum_{t=2}^{5} (y(t) - ay(t-1))^2$$

After some easy computation we find that $J_5(a) = \frac{1}{4} \left[\frac{3}{2} a^2 - \frac{3}{4} a + \frac{21}{16} \right]$ which is the measure to be **minimized**:

$$\frac{\partial J_5(a)}{\partial a} = \frac{1}{4} \left[3a - \frac{3}{4}\right] = 0$$

Which has solution

$$\hat{a}_5 = \frac{1}{4}$$

So best model identified in AR(1) class is :

$$y(t) = \frac{1}{1 - \frac{1}{4}Z^{-1}}e(t)$$

2. Compute prediction

$$\hat{y}(6|5) = \frac{1}{4}y(5) = \frac{1}{16}$$

3. First remark

Is $\hat{a} = \frac{1}{4}$ fine? Yes , it hold the condition that |a| < 1.

What if $\hat{a} = 4$? In that case we have to find the **canonical form** of the system.

4. Second remark

The variance of the white noise λ^2 is not required for the model and prediction estimation. If we wanted to estimate also λ^2 and not only a (which still is the most important parameter to estimate):

$$\lambda^2 = var[e(t)] \to var[\epsilon(t)]$$

Using an approximate estimation using the 5 available data points :

$$\hat{\lambda}_5^2 = J_5(\hat{a}_5)$$

4.2 Identification of ARMAX models

While ARX model estimation relies on *least squared method*, ARMAX model estimation relies on **maximum likelihood**.

Given the available N-long data vector ${\bf u}$ and ${\bf y}$, the model class is $m(\theta):y(t)=\frac{B(Z)}{A(Z)}u(t-1)+\frac{C(Z)}{A(Z)}e(t)$ where

$$A(Z) = 1 + \dots + a_m Z^{-m}$$

$$B(Z) = b_0 + \dots + b_p Z^{-p}$$

$$C(Z) = 1 + \dots + c_n Z^{-n}$$

and working under the hypothesis that $\frac{C(Z)}{A(Z)}$ is in **canonical representation** \rightarrow **ARMAX**(m,n,p+1) where k=1 is **not** a restriction but indeed the **most general case**.

Using the P.E.M. approach:

$$\hat{\theta} = argmin_{\theta} = \{J_N(\theta)\}$$

with performance index

$$J_N(\theta) = \frac{1}{N} \sum_{t=1}^{N} (y(t) - \hat{y}(t|t-1,\theta))^2$$

where the parameter vector has dimension $n_{\theta} = m + n + p + 1$:

$$\theta = [a_1, ..., a_m, b_0, ...b_p, c_1, ..., c_m]^T$$

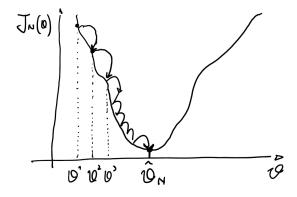
Since k=1 $\epsilon(t) = e(t) = \frac{A(Z)}{C(Z)}y(t) - \frac{B(Z)}{C(Z)}u(t-1)$ so the performance index

$$J_N(\theta) = \frac{1}{N} \sum_{t=1}^{N} \left(\frac{A(Z)}{c(Z)} y(t) - \frac{B(Z)}{C(Z)} u(t-1) \right)^2$$

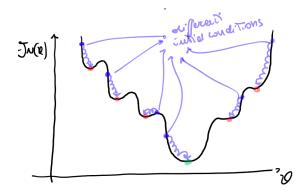
This performance index is now an issue as it kills the optimisation : C(Z) is a **non-linear** function of θ so it is not a quadratic function of θ . So the **minimization** requires an **iterative approach**.

4.2.1 Loss function: Maximum Likelihood Method

The **iterative** loss function of ARMAX models starts from an initial condition $\theta^1 \to \theta^2 \dots \to \hat{\theta}_N$ until a final solution is reached



Unfortunately usually $J_N(\theta)$ has lots of **local minima**:

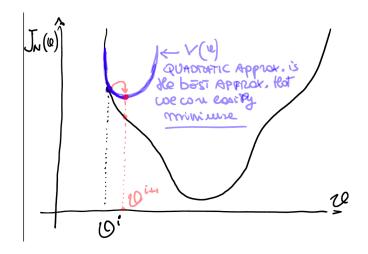


To attenuate this disturbance many **different** initial conditions θ^1 should me used, each time trying to reach the best solution. The problem is that there is no **theoretical guarantee** that solution converges to the global minimum.

The key problem for iterative methods is building the **step**:

$$\theta^i \to \theta^{i+1}$$

. The basic idea is to compute the local quadratic approximation of $J_N(\theta)$ around the point θ^i and minimize the local approximation using the Newton method



Using **Taylor's expansion** around θ^i :

$$\gamma(\theta) = J_N(\theta^i) + \left\lceil \frac{\partial J_N(\theta)}{\partial \theta} \right|_{\theta^i} \left\lceil (\theta - \theta^i) + \frac{1}{2} (\theta - \theta^i)^T \left\lceil \frac{\partial^2 J_N(\theta)}{\partial^2 \theta} \right\rceil_{\theta^i} \right\rceil (\theta - \theta^i) \dots$$

Minimizing the function $\gamma(\theta) \to \frac{\partial \gamma(\theta)}{\partial \theta} = 0$:

$$\left[\frac{\partial J_N(\theta)}{\partial \theta} \bigg|_{\theta^i} \right] + \frac{1}{2} \cdot 2 \left[\frac{\partial^2 J_N(\theta)}{\partial^2 \theta} \bigg|_{\theta^i} \right] (\theta - \theta^i) = 0$$

$$\theta = \theta^{i} - \left[\left. \frac{\partial^{2} J_{N}(\theta)}{\partial^{2} \theta} \right|_{\theta^{i}} \right]^{-1} \cdot \left[\left. \frac{\partial J_{N}(\theta)}{\partial \theta} \right|_{\theta^{i}} \right]$$

Where the first is matrix is the inverse of the second order derivative (**Hessian Matrix**) of $J_N(\theta)$ around θ^i and the second matrix is the first order (**Gradient Vector**) of $J_N(\theta)$ around θ^i .

• Gradient vector

$$\frac{\partial J_N(\theta)}{\partial \theta} = \frac{2}{N} \sum_{t=1}^{N} \epsilon(t) \cdot \frac{\partial \epsilon(t)}{\partial \theta}$$

• Hessian matrix

$$\frac{\partial^2 J_N(\theta)}{\partial^2 \theta} = \frac{2}{N} \sum_{t=1}^N \frac{\partial \epsilon(t)}{\partial \theta} \cdot \frac{\partial \epsilon(t)}{\partial \theta}^T + \frac{2}{N} \sum_{t=1}^N \epsilon(t) \frac{\partial^2 \epsilon(t)}{\partial \theta^2}$$

Which can be approximated to:

$$\frac{\partial^2 J_N(\theta)}{\partial^2 \theta} = \frac{2}{N} \sum_{t=1}^N \frac{\partial \epsilon(t)}{\partial \theta} \cdot \frac{\partial \epsilon(t)}{\partial \theta}^T$$

The approximation holds because of three reasons:

1. Reason

We can compute the Hessian using only $\frac{\partial \epsilon(t)}{\partial \theta}$ without the burden of computing $\frac{\partial^2 \epsilon(t)}{\partial \theta^2}$

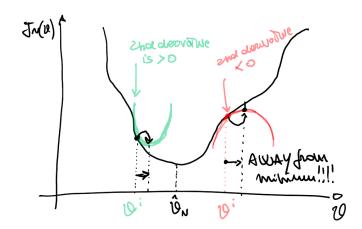
2. Reason

Since $\epsilon(t,\theta) = y(t) - \hat{y}(t|t-1;\theta)$ notice that $\frac{\partial^2 \epsilon(t)}{\partial \theta^2}$ does **not** depend on y(t) but only on $\hat{y}(t|t-1;\theta) \to y(t-1), y(t-2), \dots$

Moreover if θ^i is close to the **optimum**, $\epsilon(t;\theta) \approx e(t)$. So under these assumptions $\epsilon(t;\theta^i)$ and $\left. \frac{\partial^2 \epsilon(t)}{\partial \theta^2} \right|_{\theta^i}$ are **uncorrelated** (orthogonal).

3. Reason

By neglecting the second Hessian term we can guarantee that the approximation is **semi-definite positive** (≥ 0)



In conclusion the approximation **guarantees** that the updating step always goes in the **right direction!**

The final updating rule is:

$$\theta^{i+1} = \theta^i - \left[\sum_{t=1}^N \frac{\partial \epsilon(t; \theta^i)}{\partial \theta} \cdot \frac{\partial \epsilon(t; \theta^i)^T}{\partial \theta} \right]^{-1} \cdot \left[\sum_{t=1}^N \epsilon(t; \theta) \cdot \frac{\partial \epsilon(t; \theta^i)}{\partial \theta} \right]$$

4.2.2 Possible updating rules

There are different **updating rules** which have all the same general expression:

$$\theta^{i+1} = \theta^i - \bigcirc \cdot \left[\frac{\partial J_N(\theta)}{\partial \theta} \right]$$

• Gradient method

$$\bigcirc = \mu$$

 μ is a scalar number called **step**. Characteristics :

- + simplest method
- + correct direction guaranteed
- slow when near to minimum
- sensitive to choice of μ (too small = slow, too big = instability)

This method is the backpropagation rule in NN.

• Newton method

 \bigcirc = Inverse of hessian matrix of performance index

• Quasi- Newton method

 \bigcirc = Inverse of \ge approximation of the Hessian

This is the one seen in 4.2.1.

Has all the positive aspects of Newton method (variable step tuned to the specific point of optimisation)but can guarantee the right direction of the step.

To avoid the singularity of matrix $\sum_{t=1}^{N} \frac{\partial \epsilon(t;\theta^{i})}{\partial \theta} \cdot \frac{\partial \epsilon(t;\theta^{i})^{T}}{\partial \theta}$ usually δI is added where δ is a small number and I the identity matrix.