227-0689-00L: System Identification HS2023

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Chapter 1

Least-Squares Estimation

Consider a model

$$Y = \Phi\theta_0 + \epsilon \tag{1.1}$$

where $Y = \begin{bmatrix} y_0 & \dots & y_{N-1} \end{bmatrix}^{\top}$ is a vector containing the measurements $\{y_0, \dots, y_{N-1}\}$, $\Phi \in \mathbb{R}^{N \times p}$ is called the regressor and $\theta \in \mathbb{R}^p$ is the model parameter to be estimated; p is the model order. The noise vector $\epsilon = \begin{bmatrix} v_0 & \dots & v_{N-1} \end{bmatrix}^{\top}$ has zero mean $\mathbb{E}\left\{\epsilon\right\} = 0$ and covariance¹ $\mathbb{E}\left\{\epsilon\epsilon^{\top}\right\} = R$, a symmetric positive definite matrix.

The maximum likelihood (ML) probability² requires us to minimize

$$(Y - \Phi\theta)^{\top} R^{-1} (Y - \Phi\theta) = ||C(Y - \Phi\theta)||_2^2$$
 (1.2)

$$\epsilon = C\epsilon$$

where C is the Cholesky decomposition of R and $e \sim \mathcal{N}(0,1)$ is i.i.d.

²What is the expectation value of eq. (1.2)? Using eq. (1.5), I think the term $Y - \Phi\theta$ should be evaluated as

$$Y - \Phi \hat{\theta}_{\mathrm{ML}} = \Phi \theta_0 + \epsilon - \Phi \theta_0 - \Phi K \Phi^{\top} R \epsilon = \left(I - \Phi K \Phi^{\top} R \right) \epsilon.$$

The average

$$\mathbb{E}\left\{\left(Y - \Phi\hat{\theta}_{\mathrm{ML}}\right)^{\top}R^{-1}(Y - \Phi\hat{\theta}_{\mathrm{ML}})\right\} = \mathbb{E}\left\{\epsilon^{\top}\left(I - \Phi K \Phi^{\top}R\right)^{\top}R^{-1}\left(I - \Phi K \Phi^{\top}R\right)\epsilon\right\}$$

but I do not know what the term $\mathbb{E}\left\{\epsilon^{\mathsf{T}}R^{-1}\epsilon\right\}$ evaluates to. I expect this to be equal to N, at least this is what happens with an uncorrelated noise with constant variance.

Note also that while the matrix vv^{\top} has rank 1, the variance matrix R for Gaussian distributed white noise is full rank because it is the sum of random vectors that span the whole space.

 $^{^{1}\}epsilon$ can be numerically constructed from

where C is the Cholesky decomposition of the symmetric positive definite matrix $C^{\top}C = R^{-1}$.

The minimum of eq. (1.2) is found by setting the gradient of the expression with respect to θ to zero, which gives the normal equation

$$\Phi^{\top} R^{-1} \Phi \hat{\theta}_{\mathrm{ML}} = \Phi^{\top} R^{-1} Y.$$

The mathematical solution is given by

$$\hat{\theta}_{\mathrm{ML}} = \left(\Phi^{\top} R^{-1} \Phi\right)^{-1} \Phi^{\top} R^{-1} Y \tag{1.3}$$

and it exists provided that $C\Phi$ has full rank³: rank $(C\Phi) = p$. When this is the case, the system is said to be *persistently excited*, see Sect.. Given the freedom to choose Φ , one must select it such that $C\Phi$ is persistently excited.

Numerically one should not form the normal equation directly because it squares the condition number of $C\Phi$ and rely either on the QR decomposition or on the SVD to solve eq. (1.2). This is taken care automatically by MATLAB when using the backslash \setminus operator

$$\hat{\theta}_{\mathrm{ML}} = (C\Phi) \backslash (CY). \tag{1.4}$$

From this, it is clear that the quantities that matter are $C\Phi$ and CY.

1.1 Bias, Covariance and MSE of the Least Squares Estimation

The linear estimator eq. (1.3) is unbiased⁴:

$$\mathbb{E}\left\{\hat{\theta}_{\mathrm{ML}}\right\} = \theta_{0}.$$

$$\hat{\theta}_{\mathrm{ML}} = K \Phi^{\top} R^{-1} (\Phi \theta_0 + \epsilon) = \theta_0 + K \Phi^{\top} R^{-1} \epsilon \tag{1.5}$$

³The usual warning holds for the rank. Instead we want to have the matrix $C\Phi$ with the smallest condition number for the estimate to be numerically stable, which is a stronger condition than full rank, to uniquely determine the best estimate $\hat{\theta}_{ML}$.

⁴Recalling that R and $\Phi^{\top}R^{-1}\Phi$ are symmetric matrices, letting $K \doteq (\Phi^{\top}R^{-1}\Phi)^{-1}$ and using eq. (1.1) into eq. (1.3), we obtain

The covariance⁵

$$\operatorname{cov}\left(\hat{\theta}_{\mathrm{ML}}\right) = \left(\Phi^{\top} R^{-1} \Phi\right)^{-1}$$

Note that *only* in the case of diagonal (e.g. uncorrelate) with all elements equal $R = \sigma^2 I_N$ the covariance reduces to $\sigma^2 (\Phi^{\top} \Phi)^{-1}$. Lastly, we consider the mean squared error

$$MSE\left(\hat{\theta}_{ML}\right) = \underbrace{\left|\left|Bias\left(\hat{\theta}_{ML}\right)\right|\right|_{2}^{2}}_{=0} + tr\left(cov\left(\hat{\theta}_{ML}\right)\right)$$

which reduces to $\sigma^2 N$ for the case of uncorrelated noise.

1.1.1 Geometric Interpretation of Least-Squares

The least squares problem

$$||b-Ax||_2$$

has the following geometric interpretation: the solution is that for which the residuals $v \doteq b - Ax$ are outside (i.e. orthogonal) of the space spanned by A. In other words, we require⁶ the scalar product $\langle Az, v \rangle$ to be zero for all z:

$$0 = \langle Az, v \rangle = (Az)^{\top} (b - Ax) = z^{\top} \left(A^{\top}b - A^{\top}Ax \right) \qquad \forall z$$
$$\rightarrow A^{\top}Ax = A^{\top}b$$

which is the normal equation.

from which $\mathbb{E}\left\{\hat{\theta}_{\mathrm{ML}}\right\} = \theta_{0}$ since $\mathbb{E}\left\{\epsilon\right\} = 0$. Moreover

$$\operatorname{cov}\left(\hat{\theta}_{\mathrm{ML}}\right) = \mathbb{E}\left\{\left(\hat{\theta}_{\mathrm{ML}} - \mathbb{E}\left\{\hat{\theta}_{\mathrm{ML}}\right\}\right) \left(\hat{\theta}_{\mathrm{ML}} - \mathbb{E}\left\{\hat{\theta}_{\mathrm{ML}}\right\}\right)^{\top}\right\} \qquad \mathbb{E}\left\{\hat{\theta}_{\mathrm{ML}}\right\} = \theta_{0}$$

$$= \mathbb{E}\left\{\left(K\Phi^{\top}R^{-1}\epsilon\right) \left(K\Phi^{\top}R^{-1}\epsilon\right)^{\top}\right\}$$

$$= \mathbb{E}\left\{K\Phi^{\top}R^{-1}\epsilon\epsilon^{\top}R^{-1}\Phi K\right\}$$

$$= K\Phi^{\top}R^{-1}\mathbb{E}\left\{\epsilon\epsilon^{\top}\right\}R^{-1}\Phi K \qquad \mathbb{E}\left\{\epsilon\epsilon^{\top}\right\} = R$$

$$= K.$$

 $^{^5}$ Is there a way to understand the form of covariance matrix without going through the calculation?

⁶I believe the proof given in class is not correct: to span the full column space of A, one has to multiply by a generic vector z; in class z = x was taken.

1.2 Random Notes: The Covariance Matrix and the Choice of the Measurement Points

These are my considerations that are not part of the lecture.

- The off-diagonal elements of the covariance matrix $cov(\theta)$ represent the correlations between the errors of the variables θ . It is therefore not justified to discard them and take θ 's standard deviations as the square root of $cov(\theta)$'s diagonal elements because one discards the correlations: the ball of probability is in general an ellipse with the axes not parallel to the variable directions.
- Given the freedom to choose the measurement points, is there a "best" way of placing them?

There are two factors that determine the covariance matrix: the choice of basis and the choice of points. The choice of basis is determined by the variables that one wants to extract: a linear transformation between one basis and the other will also transform the covariance matrix and the only concern may be the numerical stability (although one should expect that the measurement errors dominate).

Determining the position of the measurement points by minimizing one (or more elements of the covariance) is in general a non-convex problem. However this is done.

Chapter 2

Regularized FIR Models

We have seen that the solution $\hat{\theta}_{\rm ML}$ of the least squares problem eq. (1.2)

$$\arg\min_{\theta}||y - \Phi\theta||_2^2$$

is unbiased, see Sect. 1.1.

We can however choose to have a biased estimate to reduce the mean square error

$$MSE\left(\hat{\theta}\right) = E||\hat{\theta} - \theta_0||_2^2.$$

This can be achieved if we modify the minimization problem by adding a regularization term

$$||Y - \Phi\theta||_2^2 + \gamma \theta^\top P^{-1}\theta \tag{2.1}$$

where γP^{-1} is a positive definite matrix. P is called the *kernel* or the regularization matrix¹.

Regularization prevents overfitting, reduces the sensitivity to noise and can improve the estimate by a proper choice of the kernel (e.g. for instance if the system is known to be stable, this information can be used to improve the estimation.)

The closed form mathematical solution to eq. (2.1) is given by

$$\hat{\theta}^{R} = \left(\Phi^{\top}\Phi + \gamma P^{-1}\right)^{-1}\Phi^{\top}Y. \tag{2.2}$$

¹When $P = \mathbb{I}_N$ the technique is called ridge regression, otherwise it goes under the name of Tikhonov regularization.

2.0.1 The James-Stein Estimator

The James-Stein estimator was the first biased estimator that had a smaller MSE compared to least squares on all the [1, page 3]. It can be cast into a ridge regression problem

$$||Y - \theta||^2 + \gamma ||\theta||^2, \qquad \gamma = \frac{(N-2)\sigma^2}{||Y||^2 - (N-2)\sigma^2}.$$

2.1 Bias, Covariance and MSE of the Regularized Least Squares Estimation

We assume uncorrelated noise² with constant variance σ^2 : $\mathbb{E}\left\{\epsilon\epsilon^{\top}\right\} = \sigma^2 I_N$. The estimate eq. (2.2) has bias, covariance³ and MSE

$$\begin{aligned} \operatorname{Bias}\left(\hat{\theta}^{\mathrm{R}}\right) &= -\left(\Phi^{\top}\Phi + \gamma P^{-1}\right)^{-1} \gamma P^{-1}\theta_{0} \\ \operatorname{cov}\left(\hat{\theta}^{\mathrm{R}}\right) &= \sigma^{2}\left(\Phi^{\top}\Phi + \gamma P^{-1}\right)^{-1} \Phi^{\top}\Phi\left(\Phi^{\top}\Phi + \gamma P^{-1}\right)^{-1} \\ \operatorname{MSE}\left(\hat{\theta}^{\mathrm{R}}\right) &= \left|\left|\operatorname{Bias}\left(\hat{\theta}^{\mathrm{R}}\right)\right|\right|_{2}^{2} + \operatorname{tr}\left(\operatorname{cov}\left(\hat{\theta}^{\mathrm{R}}\right)\right) \end{aligned}$$

The MSE is minimized by this choice of parameters⁴

$$\gamma^* = \sigma^2, \quad P^* = \theta_0 \; \theta_0^\top. \tag{2.3}$$

The optimal regularization matrix P^* is unknown because it depends on the unknown θ_0 but the approximate knowledge of the solution helps to construct a "good" P. Compared to the ridge regression which only tries to decrease $||\theta||^2$, the Tikhonov regularization can use information about the system.

$$\hat{\theta}^{\mathrm{R}} - \mathbb{E} \left\{ \hat{\theta}^{\mathrm{R}} \right\} = K \Phi^{\top} \epsilon$$

and

$$\operatorname{cov}\left(\hat{\theta}^{\mathrm{R}}\right) = \mathbb{E}\left\{\left(K\Phi^{\top}\epsilon\right)\left(K\Phi^{\top}\epsilon\right)^{\top}\right\} = K\Phi^{\top}\mathbb{E}\left\{\epsilon\epsilon^{\top}\right\}\Phi K = \sigma^{2}K\Phi^{\top}\Phi K.$$

²In Sect. 1 we considered the more general case of correlated noise $\mathbb{E}\left\{\epsilon\epsilon^{\top}\right\} = R$.

³Let $K = \left(\Phi^{\top}\Phi + \gamma P^{-1}\right)^{-1}$ so that $\hat{\theta}^{R} = K\Phi^{\top}Y$. Plugging $Y = \Phi\theta_0 + \epsilon$, we have

⁴Even for the $Y = \theta + \epsilon$ and an impulse response, so that Φ has only ones on the main diagonal (but it is not the identity matrix because in general Φ is a tall matrix), the forms of the bias and covariance for the optimal parameters do not seem to have a special form.

2.1.1 Estimation Bias

As far as we (as students) understand, there are two sources of bias:

• for an FIR system, one can underestimate the length of the response by choosing the model order too small, $p < \tau_{\text{max}}$: this is the truncation error of Sect..

For least squares we said the estimation is unbiased: this is the case only when the order is correctly selected, otherwise it is too biased (see slide 10–10). Indeed we are trying to fit the data as

$$\hat{\theta} = \arg\min_{\theta} ||Y - \Phi_p \theta|| = \left(\Phi_p^{\top} \Phi_p\right)^{-1} \Phi_p^{\top} Y$$

but the model is $Y = \Phi \theta_0 + \epsilon$. The bias then becomes

$$\mathbb{E}\left\{\hat{\theta}\right\} = \left(\Phi_p^{\top} \Phi_p\right)^{-1} \Phi_p^{\top} \Phi \theta_0 \neq \theta_0.$$

(This means the trade-off bias-variance can also be done with LS by varying the order parameter, but to a lower extent than using regularization where one has extra parameters available.)

• With the order parameter correct, one can bias the estimator on purpose by setting it to a value different from θ_0 .

2.1.2 Choice of the Regularization Matrices

A commonly used regularization matrix is the diagonally-correlated (DC) kernel

$$[P]_{ij} = c\alpha^{\frac{i+j}{2}}\rho^{|i-j|}$$

with $0 \le c, \ 0 \le \alpha \le 1$ and $-1 \le \rho \le 1$: ρ describes correlations α and the decays.

A simpler version of it is the tuned-correlated (TC) kernel where one lets $\rho = \sqrt{\alpha}$ in the DC kernel

$$[P]_{ij} = c\alpha^{\max\{i,j\}} = \begin{bmatrix} \alpha & \alpha^2 & \alpha^3 & \dots & \alpha^n \\ \alpha^2 & \alpha^2 & \alpha^3 & \dots & \alpha^n \\ \alpha^3 & \alpha^3 & \alpha^3 & \dots & \alpha^n \\ \dots & & & & \\ \alpha^n & \alpha^n & \alpha^n & \dots & \alpha^n \end{bmatrix}$$

The top left entries are large, the bottom right small. Since the inverse enters the minimization problem (cost function), the top left part is small whereas the bottom down are large and induce a larger cost.

The inverse of the DC kernel is a tridiagonal with the explicit form

$$[P^{-1}]_{ij} = \frac{c_{ij}}{1 - \rho^2} (-1)^{i+j} \alpha^{-\frac{i+j}{2}} \rho^{|i-j|}, \qquad c_{ij} = \begin{cases} 1 + \rho^2 & \text{if } i = j = 2, \dots, p-1 \\ 0 & \text{if } |i-j| > 1 \\ 1 & \text{otherwise} \end{cases}$$

During the exercises Jared commented that a TC/DC kernel should not be used for a FIR. The question is still unanswered in Moodle: https://moodle-app2.let.ethz.ch/mod/forum/discuss.php?d=137879.

Despite the TC having entries that are exponential, the TC kernel is not $\theta_0 \; \theta_0^{\mathsf{T}}$ of eq. (2.3) with

$$\theta_0 = \begin{bmatrix} a & a^2 & \dots \end{bmatrix}$$

from the truncation of a simple IIR filter $\frac{z^{-1}}{1-az^{-1}}$ for two reasons:

- $\theta_0 \, \theta_0^{\top}$ has equal powers on the diagonals parallel to the anti-diagonal;
- $\theta_0 \; \theta_0^{\top}$ has rank 1 and is therefore not invertible.

2.1.3 Why the LS May Perform Badly

Least-squares performs badly when the regression matrix Φ is ill-conditioned⁵ as a result of not being persistently excited and a small perturbation of the measurement data Y can have a large impact on the estimate $\hat{\theta}_{\text{ML}}$. Regularization helps in these cases. (Is there an easy way to show that $\Phi^{\top}\Phi + \gamma P^{-1}$ has a larger minimum singular value?)

2.2 Cross-Validation Methods

They are methods to select the discrete model orders (in least-squares or ARX) or the regularization parameters (e.g. γ and the kernel parameter α in regularized least squares).

$$\operatorname{cond}(A) = \frac{\sigma_1}{\sigma_n}.$$

Note that from a numerical point of view, the matrix being full rank is not a guarantee that the LS problem can be stably solved.

⁵The condition number of a matrix is defined as the ratio between the largest σ_1 and the smallest σ_n singular value of the matrix

A popular method is called hold-out cross-validation and consists of the following steps:

- 1. split the data into two (equal but not necessarily) non-overlapping parts, one for estimation and the other for validation. The way in which the data is split depends on the type of data: for data generated by dynamical systems, where the data is sequential in time, the order is important. Here the first $N_{\rm id}$ pairs of control inputs and output values $\{u(k), y(k)\}$ are used for identification and the remaining for validation. For static data, the order is not important: one could choose the odd-indexed pairs for identification and the even-indexed pairs for validation;
- 2. use the estimation data to estimate a model for each mode structure or each regularization parameter. In the case of the continuous parameter γ , one estimates the model on a set of gridded values. If there are two continuous parameters, e.g. γ and α , the gridding is done for both parameters;
- 3. select the model structure / regularization parameters that give a model with least prediction error on validation data;
- 4. use the selected model structure / regularization parameters and the complete data set to estimate a final model.

2.3 Numerical Solution of the Tikhonov Regularization Problem

For the numerical solution I am aware of two methods:

1. Solve directly the normal equation

$$\left(\Phi^{\top}\Phi + \gamma P^{-1}\right)\hat{\theta} = \Phi^{\top}y \tag{2.4}$$

using MATLAB's backslash operator: although this squares Φ 's condition number when constructing $\Phi^{\top}\Phi$, at least it does not require to explicitly construct the inverse;

2. Use the generalized SVD decomposition. We first manipulate eq. (2.1) to obtain

$$||y - \Phi \theta||_2^2 + ||D\theta||_2^2$$

by letting D be the Cholesky decomposition of γP^{-1} : that is $D^{\top}D = \gamma P^{-1}$.

There exists different definitions of generalized SVD: MATLAB implements the call $[U,V,X,C,S] = gsvd(\Phi, D)$ where U and V are unitary matrices, and C and S are nonnegative diagonal matrices such that

$$\Phi = UCX^{\top} \qquad D = VSX^{\top} \qquad C^2 + S^2 = I.$$

Inserting the decompositions into eq. (2.4) gives

$$X^{\top}\theta = (UC)^{\top}y$$

which can be solved as $\theta = X^{\top} \setminus (UC)^{\top} y$;

3. Eq. (2.1) can also be rewritten as

$$\left\| \begin{bmatrix} \Phi \\ D \end{bmatrix} \theta - \begin{bmatrix} y \\ 0 \end{bmatrix} \right\|_{2}^{2}$$

and solved by backslash.

I tested the three approaches on prob. 8 and found the fastest to be the third which does not construct the normal equation and is therefore better conditioned at the small cost of increasing the problem size by the r rows of P^{-1} . The first method is twice slower and the second is 100 (!) times slower: the problem is most likely the call to gsvd. This is strange because this is the method normally suggested, but probably only for larger regularization matrices.

Chapter 3

Error Prediction Methods and Transfer Function Models

We assume that the complete system model is given by

$$y(k) = G(z)u(k) + H(z)e(k).$$
 (3.1)

Prediction error-based identification methods estimate the transfer functions G(z) and H(z) by minimizing the objective function $J(\epsilon)$, a function of the prediction error ϵ

$$\hat{\epsilon}(k) = y(k) - \hat{y}(k|k-1).$$

The prediction $\hat{y}(k|k-1)$ for the time k is a function of the previous measurements k-1, k-2,... only.

3.1 Prediction

We assume that H(z) is monic¹ and stable². In this case, given the filtered noise v(k) = H(z)e(k), e(k) can be reconstructed from v(k) as

$$e(k) = H^{-1}(z)v(k).$$
 (3.2)

$$H(q) = 1 + \sum_{k=1}^{\infty} h(k)q^{-k}.$$

¹Monic means that h(0) = 1:

 $^{^{2}}H(z)$ has only poles strictly inside the unit circle.

We now seek to predict v(k) given the past values up to time k-1: this is called the *one-step ahead* estimate. Since H(z) is monic, we split the filtered noise contribution into the term e(k) and other terms up to time k-1

$$v(k) = H(z)e(k) = e(k) + (H(z) - 1)e(k)$$
(3.3)

The predicted filtered noise $\hat{v}(k|k-1)$ becomes

$$\hat{v}(k|k-1) = (H(q)-1)e(k)$$

This can be intuitively understood because the error probability function distribution for $\{e\}$ has zero mean³: if we were left to guess e(k), we would guess e(k) = 0. Making use of eq. (3.2), the one-step ahead estimate

$$\hat{v}(k|k-1) = (1 - H^{-1}(z)) v(k) \tag{3.4}$$

is determined only from the knowledge of the past values of v(k-i) with i>0.

For the model of eq. (3.1), the one-step ahead predictor

$$\hat{y}(k|k-1) = G(z)u(k) + \hat{v}(k|k-1)$$

can be rewritten with the help of the expression for $\hat{v}(k|k-1)$ as

$$\hat{y}(k|k-1) = H^{-1}(z)G(z)u(k) + (1 - H^{-1}(z))y(k). \tag{3.5}$$

The prediction error⁴

$$\epsilon(k) = y(k) - \hat{y}(k|k-1) = H^{-1}(z)v(k) = e(k)$$
(3.6)

is the noise e(k) that cannot be predicted: the *innovation* is the part of the output prediction that cannot be estimated from past measurements.

$$\hat{v}(k|k-1) = \arg\max_{x} f_e(x - m(k-1)), \qquad m(k-1) = (H(q) - 1)e(k).$$

 4 Using eq. (3.5)

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

$$= y(k) - H^{-1}(z)G(z)u(k) - (1 - H^{-1}(z))y(k)$$

$$= H^{-1}(z)(y(k) - G(z)u(k))$$

$$= H^{-1}(z)v(k) = e(k)$$

³Had the probability distribution $f_e(x)$ not had a zero mean, we would have to modify the predition according to

3.1.1 Example: Moving Average

The model

$$v(k) = e(k) + ce(k-1) \to H(z) = 1 + cz^{-1}$$

is invertible when |c| < 1. The one-step ahead predictor eq. (3.4) can be expressed in terms of the error $\epsilon(k-1)$ using eq. (3.6)

$$\hat{v}(k|k-1) = (1 - H^{-1}(z)) H(z)e(k) = cz^{-1}\epsilon(k) = c\epsilon(k-1)$$

3.2 Family of Transfer-Function Models

The advantage of the transfer-function models is that they can be described by fewer parameters and that the inputs required to identify the system do not have to be persistently exciting as it is the case when one wants to identify frequency or time-response: we have seen in Sect. ?? that for frequency domain methods, the order of excitation must be double the number of complex estimates of the transfer function $G(e^{j\omega_n})$ since gain and phase must be determined for each frequency; for a time response one requires the same persistency order as the number of impulse response terms.

If we control the input, this requirement is easy to satisfy. If on the other hand the data is given, this may not be the case and in these situations, one is better off looking for transfer functions/state space representations because of the reduced numbers of parameters to identify.

Prediction error-based identification methods construct the prediction error

$$\epsilon(k,\theta) = y(k) - \hat{y}(k,\theta)$$

from the (one-step ahead) predictor $\hat{y}(k,\theta)$ which is based on the guesses $\hat{G}(z) = G(z,\theta)$ and $\hat{H}^{-1}(z) = H^{-1}(z,\theta)$, the guesses being parametrized by θ . The optimal θ^* is the argument that minimizes the cost function $J = J(\epsilon)$

$$\theta^* = \arg\min_{\theta} J(\epsilon(k, \theta)).$$

Typical choices for the objective functions $J(\epsilon)$ are the 2-norm $||\epsilon||_2^2$ or the maximum deviation, the ∞ -norm $||\epsilon||_{\infty}$. The kind of minimization depends on how the models G(q) and $H^{-1}(q)$ are parametrised. In general, the parametrization will not be linear and the optimization may not be convex. Note moreover that the minimization of $||\epsilon||_2^2$ is not equivalent to the least squres method, unless H(z)=1:

$$y(k) = G(z)u(k) + e(k).$$

Since this approach does not require an a-priori knowlege of the system, it is also called *black-box* approach.

3.2.1 Equation Error Model Structure (ARX)

The ARX model⁵

$$y(k) = B(z)u(k) + (1 - A(z))y(k) + e(k)$$
(3.7)

is a simple input-output relationship where the error enters as a direct term. We take

$$A(z) = 1 + a_1 z^{-1} + \dots + a_n z^{-n}, \qquad B(z) = b_1 z^{-1} + \dots + b_m z^{-m}$$

note that A(z) is monic and B(z) does not have a constant term, *i.e.* the model has no feed-through. It corresponds to the model of eq. (3.1) with

$$G(z) = \frac{B(z)}{A(z)}, \qquad H(z) = \frac{1}{A(z)}.$$

and generates the one-step ahead predictor

$$\hat{y}(k|k-1) = (1 - A(z))y(k) + B(z)u(k)$$
(3.8)

either by plugging G(z) and $H^{-1}(z)$ in eq. (3.5) or by using the result of eq. (3.6) that $\hat{y}(k|k-1) = y(k) - e(k)$ in eq. (3.7).

While the model covers a limited set of real-world problems (notably those where the perturbation act as a force), the predictor eq. (3.8) has the advantage of forming the linear regressor

$$y(k) = \varphi^{\top}(k)\theta$$

the parameter vector θ being the unknown coefficients of the polynomials A(z) and B(z)

$$\theta = \begin{bmatrix} a_1 & \dots & a_n & b_1 & \dots & b_m \end{bmatrix}^\top$$

and the regressor vector $\varphi(k)$ being the output and input terms

$$\varphi^{\top}(k) = \begin{bmatrix} -y(k-1) & \dots & -y(k-n) & u(k-1) & \dots & u(k-m) \end{bmatrix}.$$

⁵From $y(k) = \frac{B(z)}{A(z)}u(k)$, rearranging the terms

$$A(z)y(k) = B(z)u(k) \to y(k) = B(z)u(k) + (1 - A(z))y(k).$$

3.2.2 Estimate Bias

The ARX model structure is applied also to systems that do not have this noise model for the simple reason that it is linear. One must be aware that this induce a bias in the estimates $G(z, \theta^*) = \frac{B(z, \theta^*)}{A(z, \theta^*)}$ and $H^{-1}(z, \theta^*)$, see the numerical example in the slides 9.30.

3.2.3 ARMAX Model Structure

The ARX model structure is not very flexible with regards to the noise model: indeed it requires the noise to have the particular structure 6 1/A(z). The ARMAX transfer function model is in the form

$$A(z)y(k) = B(z)u(k) + C(z)e(k)$$
(3.9)

where A(z) and B(z) are as in ARX and C(z) is monic. It corresponds to the model of eq. (3.1) with

$$G(z) = \frac{B(z)}{A(z)}, \qquad H(z) = \frac{C(z)}{A(z)}$$

and with the one-step ahead predictor⁷

$$\hat{y}(k|\theta) = (1 - A(z))y(k) + B(z)u(k) + (C(z) - 1)\underbrace{(y(k) - \hat{y}(k|\theta))}_{\doteq \epsilon(k,\theta)}.$$
(3.10)

Introducing the regression vector

$$\varphi^{\top}(k,\theta) = \begin{bmatrix} -y(k-1) & \dots & u(k-1) & \dots & \epsilon(k-1,\theta) & \dots \end{bmatrix}$$

eq. (3.10) induces the pseudolinear regression

$$\hat{y}(k,\theta) = \varphi^{\top}(k,\theta)\theta.$$

Working out the expression directly is cumbersome: by using eq. (3.6) in the model eq. (3.9), we have

$$\hat{y}(k|k-1) = y(k) - e(k) = \frac{B(z)}{A(z)}u(k) + \left(\frac{C(z)}{A(z)} - 1\right)e(k).$$

By plugging in the expression for the noise

$$e(k) = C^{-1}(z)(A(z)y(k)) - B(z)u(k)$$

we obtain

$$C(z)\hat{y}(k|k-1) = B(z)u(k) + (C(z) - A(z))y(k).$$

Some manipulation is still required to bring it into the final form eq. (3.10).

⁶It has the advantage (together with FIR models) of forming a linear regressor.

⁷The one-step predictor is more easily obtained by plugging the expressions for G(z) and $H^{-1}(z)$ into eq. (3.5).

3.2.4 Constrained Minimization

The optimization-based algorithm

$$\min_{\theta,\epsilon} ||\epsilon||_2$$
 subject to $Y = \Phi(\epsilon)\theta + \epsilon$

has a non-linear dependence on θ in the affine constrained.

3.2.5 General Family of Model Structures

The most general family of model structure is [2, Sect. 4]

$$A(z)y(k) = \frac{B(z)}{F(z)}u(k) + \frac{C(z)}{D(z)}e(k)$$
 (3.11)

Some of the common cases that we have seen so far are summarized in Table 3.1.

Polynomials	Name of Model Structure
В	FIR
АВ	ARX
АВС	ARMAX
A B C D	ARARMAX
B F C D	Box-Jenkins

Table 3.1: Some common black-box SISO models using the polynomials of eq. (3.11).

Bibliography

- [1] Regularized System Identification: Learning Dynamic Models from Data. Springer, 2021.
- [2] System Indentification: Theory for the User. Second edition. Prentice Hall PTR, 2009.