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Initializing Wiener-Hammerstein Models Based on Partitioning of the Best Linear Approximation

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Abstract

This paper describes a new algorithm for initializing and estimating Wiener-Hammerstein models which consist of two linear parts with a static nonlinearity in between. The algorithm makes use of the best linear model of the system, which is a consistent estimate of the systems dynamics for Gaussian excitations. The linear model is split in all possible ways into two sub-models. For all possible splits, a Wiener-Hammerstein model is initialized which means that a nonlinearity is introduced in between the two sub-models. The linear parameters of this nonlinearity can be estimated using least-squares. All initialized models can then be ranked depending on the fit. Typically, one is only interested in the best one, for which all parameters are fitted using prediction error minimization.

The paper explains the algorithm in detail and consistency of the initialization is proven. Computational aspects are investigated, showing that in most realistic cases, the number of splits of the initial linear model remains low enough to make the algorithm useful. The algorithm is illustrated on an example where it is shown that the initialization is a tool to avoid many local minima.

Key words: Wiener-Hammerstein systems, Hammerstein systems, Wiener systems, nonlinear system identification

1 Introduction

There has always been a need to identify nonlinear systems using measured data. In real life all systems are nonlinear to some extent but linear system theory and linear system identification methods have often been successfully applied. Theory for linear system identification is a fairly mature area, well covered in books like, eg. [11, 19], focusing on time domain methods and [12], focusing on frequency domain methods.

During the last two decades there has been a growing interest to go beyond linear identification and to identify real nonlinear models of the nonlinear systems. General model structures, like Volterra expansions, or different types of neural network expansions which can approximate “almost” any nonlinear relations have been proposed and investigated extensively in the literature, see, eg. [18] for an overview. The drawback using very flexible models is an increased risk of a high variance contribution to the model error and problems with the “curse of dimensionality”. The papers [8, 18] explains these problems and presents some remedies.

One approach to avoid the problems with the flexible, general, model structures is to work with block-oriented nonlinear models built up by combining blocks which are either linear dynamic or static nonlinear. By constraining the dynamics to be linear one can also use part of the linear system theory and analysis of the model properties become easier. The simplest block-oriented models are the Hammerstein model and the Wiener model where the Hammerstein model consists of a static nonlinearity followed by a linear dynamic block. The Wiener model has the two blocks in the reversed order. These models can be generalized into Hammerstein-Wiener models, which has two static nonlinear blocks with a linear dynamic block in between, and Wiener-Hammerstein models with two linear blocks with a static nonlinear block in the middle, see Figure 1.

With the introduction of a block structure, the need, and possibility, for tailored estimation methods follows. This papers presents a novel consistent algorithm for es-
estimating Wiener-Hammerstein models. In the literature, many approaches for estimating block-oriented models have been published. Early results can be found in [1] and [3]. The first one of these also contains many references to pioneer works. A stochastic embedding of the estimation problem is given in [4] and the maximum likelihood estimate is formulated. Recursive estimation of MIMO Wiener-Hammerstein models is investigated in [2].

A key problem in identifying Wiener-Hammerstein models is the initial estimates of the two linear blocks, including the degree of them. While initialization and estimation methods exists for the simpler block models, see e.g [16, 20–22] for Wiener models, and [7] for an overview of Hammerstein models, there is a lack for of these for general Wiener-Hammerstein models. In [5] an iterative initialization procedure is proposed which requires specially designed periodic excited input signals. This experimental requirement is relaxed in [15].

In [23] random, stable, initialization of the two sub-systems is investigated with quite good results. The authors show that there are many local minima, so, typically, the estimation needs to be repeated several times with different starting values to increase the chances of finding a model corresponding to a good local minimum.

In this paper a novel algorithm is proposed which is based on the Best Linear Approximation (BLA) of the nonlinear system. That means that the algorithm starts with linear identification to obtain the BLA model. Under some mild assumptions, see Section 4, the BLA model will be a consistent estimate of the concatenated dynamics of the two linear sub-systems. The core idea of the algorithm is then a brute-force approach, simply splitting the BLA model into two sub-models in all possible ways, initialize a Wiener-Hammerstein model with each one of these splits, and then choose the best initialization. The contributions of the paper are the consistency of the algorithm and that the initializations can be formulated as least squares (LS) problems which makes it feasible up to model orders of approximately 10. The consistency result means that problems with local minima are avoided which cannot be guaranteed by any other suggested algorithm for this model structure.

The idea to base the initial Wiener-Hammerstein model on the BLA model is not new. In, for example, [9] an algorithm is given where initially the entire BLA model is put on both sides of the nonlinearity.

The paper is organized as follows. Section 2 details the problem formulation and the proposed algorithm is presented in Section 3. The consistency of the algorithm is shown in Section 4. Sections 5 and 6 concern the complexity and the computational burden for the algorithm. In Section 7 the algorithm is illustrated on an example and the paper ends with conclusion in Section 8.

2 Problem Formulation

The problem formulation is divided into three steps, the definition of the model structure, the assumptions on the data, and the definition and computation of the estimate.

2.1 Model structure

The concerned model structure is of Wiener-Hammerstein type described by

\[
\begin{align*}
    z(t) &= G_1(q^{-1}, \alpha) u(t) \\
    x(t) &= f(\beta, z(t)) \\
    \hat{y}(t, \theta) &= G_2(q^{-1}, \gamma) x(t)
\end{align*}
\]

(1)

where \( \hat{y}(t) \) is the model’s prediction of the output \( y(t) \), and \( G_1(q^{-1}, \alpha) \) and \( G_2(q^{-1}, \gamma) \) are linear time invariant transfer functions in the delay operator \( q^{-1} \), and parameterized with \( \alpha \) and \( \gamma \), respectively. In the sequel, the arguments of the transfer functions are sometimes omitted to simplify the notation. The function \( f \) is a static nonlinearity parameterized with \( \beta \). All parameters of the model structure are stored in a common parameter vector

\[
\theta = [\alpha, \beta, \gamma].
\]

The first linear part of the model can be described as

\[
G_1(q^{-1}, \alpha) = \frac{b_0^1 + b_1^1 q^{-1} + \ldots + b_{m_1}^1 q^{-m_1}}{1 + a_1^1 q^{-1} + \ldots + a_{m_2}^1 q^{-m_2}}
\]

(3)

where \( \alpha = [b_0^1, \ldots, b_{m_1}^1, a_1^1, \ldots, a_{m_2}^1] \), and \( G_2(q^{-1}, \gamma) \) is described similarly with \( \gamma = [b_0^2, \ldots, b_{m_2}^2, a_1^2, \ldots, a_{m_2}^2] \).

The static nonlinearity is described as a basis function expansion

\[
f(\beta, z) = \sum_{k=1}^{n} \beta_k^1 f_k(\beta_k^2, z)
\]

(4)

where \( f_k \) are basis functions, and \( \beta \) has been divided into \( \beta^1 \), which enters linearly in \( f \), and \( \beta^2 \), which can contain several parameters, and which enters non-linearly in \( f \).

With this general description of the static nonlinearity most specific basis function expansions can be described.
with a specific choice of the basis \( f_k \). If, for example, a polynomial model is chosen, then

\[
f(\beta, z) = \beta_0^1 + \beta_1^1 z + \beta_2^2 z^2 + \ldots \beta_n^n z^n
\]

and in this case there are no parameters in \( \beta^2 \).

To define a model in this model structure, not only the parameters need to be determined but also the orders of the sub-models, and the type of basis function expansion in \( f \).

2.2 Data

For the estimation of the parameters in the model (1) a data set is assumed to be available, \( \{u(t), y(t)\}_{t=1}^N \) of \( N \) input \( u(t) \) and output \( y(t) \) samples.

The data generating process does not need to be of the form (1), but when consistency is studied in Section 4 it is assumed to be so, with a slightly more general assumption on the nonlinearity. In that case, it is important to notice that the intermediate variables \( z(t) \) and \( x(t) \) are not available. Also, for the consistency it is assumed that the input signal \( u(t) \) is Gaussian.

Otherwise rather mild conditions are needed on the data like the one in [10] for the algorithm to be applicable.

2.3 Computing Estimation

A standard prediction error approach is assumed to be used to define the estimate \( \hat{\theta}_N \) of the parameter vector \( \theta \) for the model (1) based on the data set \( \{u(t), y(t)\}_{t=1}^N \).

It is based on minimizing the prediction error

\[
\varepsilon(t, \theta) = y(t) - \hat{y}(t, \theta), \quad (5)
\]

the difference between the measured output \( y(t) \) and the prediction with (1),

\[
\hat{y}(t, \theta) = G_2(q^{-1}, \gamma) f(\beta, G_1(q^{-1}, \alpha) u(t)). \quad (6)
\]

This is done by using a criterion of fit

\[
V_N(\theta) = \frac{1}{N} \sum_{t=1}^{N} \varepsilon^2(t, \theta) \quad (7)
\]

and then defining the estimate as

\[
\hat{\theta}_N = \arg \min_{\theta} V_N(\theta). \quad (8)
\]

The estimate (8) is computed with a gradient based iterative algorithm. Given a start value \( \theta^{(0)} \), iterate

\[
\theta^{(i+1)} = \theta^{(i)} - R_i \frac{dV_N(\theta)}{d\theta} \quad (9)
\]

until convergence. The matrix \( R_i \) is to modify the search direction and step size to assure downhill steps. Depending on how \( R_i \) is chosen (9) describes a wide class of well-known standard algorithms like the Gauss-Newton and the Levenberg-Marquardt algorithms. Also the algorithm used in [23] can be obtained by choosing \( R_i \) to be the Gauss-Newton approximation with some of the smallest eigenvalues truncated.

All three blocks of the model structure (1) contain a gain parameter, and two of them are typically fixed in the iterative minimization, eg, \( b_0^1 \) and \( b_0^2 \).

In the example in Section 7 a Levenberg-Marquardt algorithm is used to compute the minimization. A software package for the Mathematica platform, [17], has been used which has the advantage that when the model structure (1) has been defined, all expressions needed in (9) to apply the Levenberg-Marquardt algorithm are calculated automatically. The symbolic features of Mathematica are used to automatically generate expressions for the derivatives. The symbolic expressions are also simplified so that, eg, multiple identical expressions are only calculated once.

Typically, \( V_N(\theta) \) can have many minima and the success of the minimization depends on the initial estimate \( \theta^{(0)} \).

The contribution of this work is a novel algorithm to compute this initial estimate.

3 Proposed Algorithm

The algorithm consists of the following steps.

Algorithm 1

1. Start with the best linear model \( G(q^{-1}) \).

Only the plant model is of interest. One can, hence, constrain to evaluate output error models. If Box-Jenkins models, or ARMAX models are used, only the estimated plant model is retained. Also frequency domain methods can be used to obtain the model. In that case a non-parametric noise weighting can be used to improve the quality of the initial estimate.

2. Split the linear model into all possible \( G_1(q^{-1}) \) and \( G_2(q^{-1}) \) so that \( G(q^{-1}) = G_1(q^{-1})G_2(q^{-1}) \).

To do this, poles and zeros of the linear model need to be calculated. These are then divided in all possible ways into sub-models \( G_1 \) and \( G_2 \). Depending on prior knowledge of the system, some of the divisions can be excluded. This is discussed in Section 5.

3. For all partitions of the linear model, \( \{G_1, G_2\} \), use \( u(t) \) and \( G_1 \) to decide values for \( \beta^2 \) and then LS to
fit the linear parameters, $\beta_1$, in the nonlinearity as initialization.

The position parameters $\beta_2$ for the basis functions are decided using the distribution of the input to the nonlinearity $\{z(t) = G_1(q^{-1})u(t)\}_{t=1}^{N}$. Minimizing (7) with respect to the parameters $\beta_1$ in (4) is straightforward by first writing (6), $\hat{y}(t, \theta)$, as

$$\hat{y}(t, \theta) = \sum_{k=1}^{n} \beta_k G_2(q^{-1})f_k(\beta_k, z(t)) = \beta^T \varphi(t)$$

where

$$\varphi^T(t) = [G_2(q^{-1})f_0(\beta_0, z(t)), \ldots, G_2(q^{-1})f_n(\beta_n, z(t))]$$

Since (10) is a linear regression, the LS estimate is given by

$$\hat{\beta}^1 = \left( \frac{1}{N} \sum_{i=1}^{N} \varphi^T(t)\varphi(t) \right)^{-1} \frac{1}{N} \sum_{i=1}^{N} \varphi^T(t)y(t).$$

4. Order the initialized models with respect to their initial fit.

This means that $V_N(\theta)$, (7) is calculated for all initializations and that the models are ranked with this measure.

5. Fit all parameters of the best, or some of the best models.

This means that the minimization algorithm (9) is applied and it is actually not part of the initialization, but the step after the initialization.

4 Consistency of the Algorithm

If the system generating the data is within the model structure described by the Wiener-Hammerstein model (1), and with some assumption on the input signal, then the consistency of the proposed algorithm when the number of data goes to infinity follows almost immediately. We make use of the following lemma.

Lemma 1 Suppose the input data $u(t)$ is a stationary normally distributed sequence and the output $y(t)$ is obtained by filtering $u(t)$ through a system of form (1), possibly corrupted by additive noise, with linear parts $G_1^0$ and $G_2^0$ being stable, single input-single output finite order transfer functions, and nonlinear part $f^0_0$, a continuous function $\mathbb{R} \rightarrow \mathbb{R}$, with limits on its amplitude and its derivative, $|f^0_0| \leq f_{\text{max}}^0$ and $|f^{0'}_0| \leq f_{\text{max}}^{0'}$. Then the best linear approximation (BLA), converges with probability 1 to

$$\kappa G_1^0(q^{-1})G_2^0(q^{-1})$$

where $\kappa$ is a constant whose value depends on $u(t)$, $f^0_0$, $G_1^0$ and $G_2^0$.

Proof See [12].

The essence of this lemma is that the BLA captures the dynamics of the two linear parts and the nonlinear function is approximated with a constant. The miss-match between the true nonlinear function and the constant is captured as noise. This is further described in [6,13,14]. There is an evident special case which might cause problems. If there is a pole in one of the linear parts which is identical to a zero in the other linear part then this zero-pole pair will be missed in the model selection step in the linear system identification of the BLA. The algorithm would then start with a BLA model with too low degree. This special case cannot be handled with this algorithm.

Given this lemma, it follows that, asymptotically in $N$, one of the partitioning of the linear models will have the correct dynamics, ie, $\kappa_1 G_1^0$ in the first linear part and $\kappa_2 G_2^0$ in the second one, where $\kappa_1 \kappa_2 = \kappa$. It remains to show consistency in the estimate of the nonlinear function. This is done using a polynomial expansion and Weierstrass’s Approximation Theorem for polynomials. The convergence is over any bounded interval.

Theorem 1 Suppose the data is generated in the same way as in Lemma 1 with the additional assumption that $G_1^0G_2^0$ does not contain any zero-pole cancellations.

Assume Algorithm 1 is applied and the nonlinearity is modeled by a polynomial expansion $f_0$ of degree $n$ combined with a saturation to limit the output for large positive and negative values outside a region which grows with the number of data. Then, on any interval $[z_a, z_b]$ for the nonlinear part of the model, the best initialized Wiener-Hammerstein model converges to the true data generating system when $N, n \rightarrow \infty$ under the constraint $n/N \rightarrow 0$.

Proof For simplicity, assume that $\kappa_1 = \kappa_2 = \kappa = 1$. The general case can be transformed so this holds. The convergence of the linear parts are already given in Lemma 1. If the linear parts are equal to the true linear parts, $G_1^0$ and $G_2^0$, then the convergence of the nonlinear function follows directly from Weierstrass’s Approximation Theorem for polynomials since the distribution of the output signal from $G_1^0$ is dense in $[z_a, z_b]$.

It remains to show the combined convergence, that the diminishing estimation error of the linear parts also gives a diminishing estimation error on the nonlinear function when $N, n \rightarrow \infty$ under the constraint $n/N \rightarrow 0$. 


The nonlinear part of the model, \( f^n(\beta, z) \), is described as
\[
f^n(\beta, z) = \begin{cases} 
\beta_0 + \beta_1 z + \beta_2 z^2 + \ldots + \beta_n z^n & |z| \leq \delta \sigma_z \\
\beta_{\max} & \text{else}
\end{cases}
\]
where \( \sigma_z \) is the standard deviation of \( z(t) = G_0^0 u(t) \) and \( \beta = [\beta_0, \beta_1, \ldots, \beta_n] \), and where \( \beta_{\max} \) is an arbitrary fixed non-negative real value which sets the nonlinearity output constant for large values of \( |z(t)| \).

Weierstrass’s Theorem guarantees that for any value of \( \delta \),
\[
\int_{-\delta \sigma_z}^{\delta \sigma_z} (f^n(z) - f^n(z))^2 dz \to 0
\]
for the minimizing values of the parameter vector \( \beta \) when \( n \to \infty \).

Now, let also \( \delta \to \infty \) but with the constraint that (14) holds so that the nonlinearity is approximated over a growing domain, of which \([z_a, z_b]\) is a subset.

The model output becomes
\[
\hat{y}(t, \theta) = G_2(f^n(G_1 u(t))) = G_2^0(f^n(G_1 u(t))) + \Delta G_2(f^n(G_1 u(t)))
\]
where the second linear model part has been described as the limit model plus the deviation from the limit model, \( \Delta G_2 \). Describing the first linear part and the nonlinear function in the same way leads to
\[
\hat{y}(t, \theta) = G_2^0(f^n(G_1 u(t))) + G_2^0(\Delta f(G_1 u(t))) + \Delta G_2(f^n(G_1 u(t)))
\]
where the last term goes to zero due to that \( \Delta G_2 \to 0 \). The middle term goes to zero if \(-\sigma_z < G_1 u(t) < \delta \sigma_z \). Since (14) and otherwise it is limited by the gain of \( G_2^0 \) times \( \int_{-\delta \sigma_z}^{\delta \sigma_z} f^n(z) \, dz \). The \( \beta \) values for which \( G_1 u(t) \) is outside the interval will be more sparse when \( \delta \) increases and that will be used below. The first term in (16) can be expanded with a Taylor expansion
\[
f^n(G_1 u(t)) = f^n(G_0^0 u(t)) + \frac{df^n}{dz}
\]
Keep only first order terms in \( \Delta G_1 \), using the fact that \( y(t) = G_2^0 f^n(G_0^0 u(t)) + \varepsilon(t) \), where \( \varepsilon(t) \) is the noise term, to form the prediction errors gives
\[
\varepsilon(t, \theta) = G_2^0 \frac{df^n}{dz} \Delta G_1 u(t) + G_2^0(\Delta f(G_1 u(t))) + \Delta G_2(f^n(G_1 u(t))) + \varepsilon(t).
\]
Squaring this and forming the criterion \( V_N(\theta) \), (7) gives the noise term which converges to the noise covariance and a number of terms which trivially can be shown to converge to zero. The only non trivial term is
\[
\frac{1}{N} \sum_{i=1}^{N} (G_2^0(\Delta f(G_1 u(t))))^2.
\]
Consider first the terms in the sum for which \( |G_1 u(t)| > \delta \sigma_z \). As described above, the \( t \) for which this holds will be a decreasing part of all time indices as \( \delta \) grows. Since the contribution of each such term is bounded, the sum of the squares of these terms will vanish when \( N \to 0 \). Remaining terms, for which \( |G_1 u(t)| \leq \delta \sigma_z \) converge to zero if \( \Delta f \to 0 \). This is possible due to (14) and since this make \( V_N(\theta) \) attain the value of the noise variance, it is the lowest possible value for \( V_N(\theta) \). Hence, \( V_N(\theta) \) is minimized when \( \Delta f \to 0 \) in the chosen interval. \( \Box \)

The theorem tells us that the proposed algorithm is sound and it gives us the true system description asymptotically. In practical situations the number of data is limited and the theorem motivates the use of the algorithm to obtain the initial parameter estimate before the iterative minimization. There are a number of comments one can make on the theorem.

- The properties of \( u(t) \) influence the convergence speed when \( N \to \infty \). Generally, an input signal which excites the system more will accelerate the convergence.
- The consistency is proven for a general continuous nonlinear function but only on an arbitrary chosen interval. A larger interval will typically make the convergence slower in \( n \) and then also in \( N \) due to the requirement \( n/N \to 0 \).
- Instead of a polynomial model for the non-linearity, any basis function expansion can be used.

5 Number of linear partitions

An obvious possible disadvantage of the algorithm is that the original linear model can be partitioned in many different ways which leads to many LS problem to solve. In this section combinatorics is used to show that for moderate orders of the BLA model the number of LS problems is reasonable.

Considering a general linear system with \( m \) poles and \( m \) zeros, one obtains,

\textbf{Result 1} A linear system of order \( m \) can be partitioned into between \( 2^m \) and \( 2^m \) different pairs of linear systems. The lower limit hold if all poles and zeros are complex and the upper limit when they are all real.

For a 10th order system this means that there will be between 1024 and \( 1024^2 > 10^6 \) possible splits of the
linear system. The lower limit is feasible for the proposed algorithm but the upper limit would typically be too many possibilities to investigate.

The number of partitions can, however, be reduced with prior knowledge or assumptions. In many situations it could make sense to require the two linear sub-systems to be proper. This gives the following result,

**Result 3** A linear system of order \( m \) can be partitioned into \( 2^{m/2+1} \) and \( 2^{m+1} \) different pairs of proper linear systems. The lower limit holds if all poles and zeros are complex and the upper limit when they are all real.

This means that for an original linear plant of order 10 the number of splits becomes between \( 2^{10/2+1} = 64 \) and \( 2^{11} = 2048 \). Note that the two special cases where all the dynamics is either placed in the first or in the second linear sub-system corresponds to Wiener and Hammerstein models, respectively. Tests of these model structures are hence treated as special cases of the Wiener-Hammerstein model selection.

Another type of possible assumption could be the order of the two linear sub-system. One then obtains,

**Result 2** A linear system of order \( m \) can be partitioned into two proper linear parts, one of order \( k \) and the other of order \( m - k \), in between \( 2^{(m/2)} \) and \( 2^{(m)} \) different ways. The lower limit is for the case that all poles and zeros are complex and the upper limit is for the case that they are all real.

For our example system of order 10, and a split into a 4th and a 6th order systems this means between \( 2^{(10/2)} = 20 \) and \( 2^{(10)} = 2 \cdot 10!/(4!6!) = 420 \) different possible splits. The number of possibilities is slightly higher if the 10th order system is to be split into two 5th order systems. However, this can only been done if at least two of the poles and two of the zeros are real. Hence the lower number of possible splits becomes \( 2 \cdot 2^{(5)} = 24 \), where the extra factor 2 is due to the two possibilities how to divide the real poles and zeros. Similarly, the upper limits of possible splits becomes \( 2^{(5)} = 2 \cdot 10!/(5!5!) = 504 \).

### 6 Computational Aspects

The main computational burden of the algorithm are in the filtering of the basis functions of the nonlinearity, (11), and in the solution of the LS problem, (12). Although each of these steps are fast, a high order linear model would lead to that these computations need be repeated many times, as described in the previous section.

The authors claim that model orders up to order 10 are feasible, at least if there are some complex poles. Of course, the exact number depends on the available computer and on how long the user is prepared to wait but since the number of linear partitions increases exponential with the model order \( m \) the feasible model order is at least “close” to 10.

To speed up computations, there are a couple of straightforward measures one can take. First, since the initializations are independent from each other, they can be executed in parallel. Most computers today have several cores and software platforms, such as Mathematica and Matlab, support parallelization without too much work for the user.

Another possible speed up is to avoid the multiple filtering, (11), of each term of the derivative of the nonlinearity through the second linear part, \( G_2(q^{-1}) \) to obtain \( \varphi(t) \). Instead, the output \( y(t) \) can be filtered through the inverse of \( G_2(q^{-1}) \) and the LS estimate is performed between the derivative of the nonlinearity and \( G_2^{-1}(q^{-1})y(t) \). This has the drawback that also the additive noise in the output will be filtered, and possibly amplified which disturbs the initial estimate. To some extent this can be compensated for using the knowledge of \( G_2(q^{-1}) \) to choose weighting filters for the LS fit.

### 7 Example

In this example the proposed algorithm is tested on data generated by a Wiener-Hammerstein model. That is, the plant is in the model set and the main message is to illustrate that the algorithm has a good chance to avoid bad local minima. The software described in [17] is used for implementing the example.

A white Gaussian signal with standard deviation 15 is used as input signal to a system with the Wiener-Hammerstein structure. Figure 2 depicts the poles and zeros of the two linear parts and Figure 5 a) shows the nonlinearity which is between them. The exact mathematical definition of the system is given in Appendix A.

![Fig. 2. Poles (x) and zeros (o) of the two linear parts of the true system.](image)

With this true system, 2000 data samples where generated, \( \{u(t), y(t)\}_{t=1}^{2000} \) where the output is corrupted with white Gaussian noise with standard deviation 0.1. This gives a signal-to-noise amplitude ratio of 57.
The first step of the algorithm is to estimate a linear model. In the example the search for the best linear model is skipped and the prior knowledge that it should be of 5th order is used. See any standard literature on system identification for strategies to obtain the best linear model, eg [11,12,19]. In Figure 3 the simulation with the 5th order linear model is shown together with the true output. The poles and zeros of the linear model are shown in Figure 4 a). Two zeros are clearly incorrect compared to the true positions shown in Figure 2 but the rest of the poles and zeros seem to be reasonably well estimated.

The linear model is now partitioned into two linear parts which should be at least first order systems. This gives 22 partitions and all of these are extended with a first order spline with 8 knots, ie, a local linear function with 9 segments. The positions of the knots are distributed so that each segment contains the same number of data.

Given the division of the linear model and the positions of the knots, the nonlinearity can be initialized according to (12). The models are sorted according to their fit at the initialization and the result is shown in Figure 4 b). The best initialization has a RMS fit of 1.05 compared to 1.75 for the linear model. The initialization of all 22 models took 35 CPU seconds on an average PC of the time of this paper. The bullets shows the RMS after initialization. The best initialization outperforms the other. In the same figure the fit after that all parameters have been fitted, is also shown with squares. The best initialized model is also the one which gives the best fit after that the criterion has been minimized with respect to all parameters. It has an RMS fit of 0.955 which corresponds with the noise level.

In Figure 5 a) the estimated nonlinearity at initialization and for the final model are shown together with the true nonlinearity.

Since poles and zeros cannot move from one linear part to the other, it is clear that the most important issue is to have the right number of poles and zero in each one of them. For this example, out of the 22 partitions, 8 of these have the correct number of poles and zeros in each linear part. That is, all of these could converge to the best model if they would not get caught in any local minima. Figure 5 b) shows the criterion decrease for these 8 models.

Clearly, the proposed initialization gives a good initialization so that only 10 iterations are needed in the recursive minimization. Some of the other initializations might lead to a good minimum if many more iterations would be applied, but after 30 iterations only a few of them have not yet terminated due to convergence to a local minimum, and their decrease is very slow.

8 Conclusions

A novel algorithm for initializing Wiener-Hammerstein models has been proposed. It starts with a best linear model which is partitioned in all possible ways into two linear sub-models. Using LS the nonlinearity between the two linear parts can efficiently be initialized. The initialized models can be ranked, and the one with the best fit is the one with best chances to converge to the global minimum when all parameters are estimated simultaneously. Consistency of the algorithm has been shown. Finally the algorithm has been illustrated with an example which shows that it can really make a difference.
Acknowledgment

This work is sponsored by the Fund for Scientific Research (FWO-Vlaanderen), the Flemish Government (Methusalem 1), and the Belgian Federal Government (IUAP VI/4). The research was performed during Jonas Sjöberg’s sabbatical at Vrije Universiteit Brussel, 2009-2010.

A True Wiener-Hammerstein system

The data used in the example in Section 7 where generated by the Wiener-Hammerstein system with of the following form of the two linear parts \( G_1(q^{-1}) = (q^2 - 1.4q + 0.5)/q^2 - 1.6q + 0.8 \) and \( G_2(q^{-1}) = (0.01q^3 - 0.014q^2 + 0.0098q)/(q^3 - 2.8q^2 + 2.6528q - 0.850944) \). The poles and zeros of these transfer functions are depicted in Figure 2. The nonlinearity, depicted in Figure 5 a), is formally defined as

\[
 f(z) = \begin{cases} 
 1.5z + 4.5 & z < -15 \\
 0.3z - 13.5 & -15 < z < -5 \\
 3z & -5 < z < 5 \\
 0.3z + 13.5 & 5 < z < 15 \\
 1.5z - 4.5 & 15 < z 
\end{cases} \quad (A.1)
\]

References


