# On the behavior of autonomous Wiener systems 

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#### Abstract

Wiener systems are nonlinear dynamical systems, consisting of a linear dynamical system and a static nonlinear system in a series connection. Existing results for analysis and identification of Wiener systems assume zero initial conditions. In this paper, we consider the response of a Wiener system to initial conditions only, i.e., we consider autonomous Wiener systems. Our main result is a proof that the behavior of an autonomous Wiener system with a polynomial nonlinearity is included in the behavior of a finite-dimensional linear system. The order of the embedding linear system is at most $\binom{n+d}{d}$ - the number of combinations with repetitions of $d$ elements out of $n$ elements - where $n$ is the order of the linear subsystem and $d$ is the degree of the nonlinearity. The relation between the eigenvalues of the embedding linear system and the linear subsystem is given by a rank-1 factorization of a symmetric $d$-way tensor. As an application of the result, we outline a procedure for exact (deterministic) identification of autonomous Wiener systems.


Key words: Block-oriented models; Wiener system; Behavioral approach, System realization, Nonlinear system identification.

## 1 Introduction

Interconnections of linear dynamic and nonlinear static systems is a popular class of nonlinear systems, refered to as block-oriented models [4,1,15]. Block-oriented models are simpler to identify from data and simpler to use for simulation and control due to the restriction of the nonlinear subsystems to be static. Among the variety of block-oriented models, the simplest special case is the Wiener system. A Wiener system consists of a linear system followed by a nonlinear static system. Despite of its limited modeling power in comparison to other block-oriented models, the Wiener system is the natural first step in the study of the class of block-oriented models and has practical applications.

A special case of an input-output system when the input dimension is zero is the autonomous system. To the best of our knowledge, currently there are no methods for autonomous Wiener system identification. Existing methods for input/output Wiener system identification [18,5,2] impose as an assumption persistency of excitation of the input. This makes them unsuitable for the autonomous case (where the input is missing or, equivalently, it is a zero sequence). In contrast, linear time-invariant identification methods such as the prediction error and subspace methods can deal seamlessly with the autonomous case.

Our main result is that an autonomous Wiener system with a polynomial nonlinearity is embedded in a finite-dimensional
linear system. In order to outline the result, consider an autonomous Wiener system $\mathscr{B}_{\mathrm{w}}$, that is a series connection of an order- $n$ linear time-invariant subsystem $\mathscr{B}$ and a degree$d$ polynomial nonlinearity $g$. We prove that $\mathscr{B}_{\mathrm{w}}$ is included in a linear time-invariant system of order $n_{\mathrm{w}} \leq\binom{ n+d}{d}$ - the number of combinations with repetitions of $d$ elements out of $n$ elements. Moreover, there is a relation between the eigenvalues of the embedding system and the eigenvalues of $\mathscr{B}$ : an eigenvalue of the embedding system is a product of up to $d$ eigenvalues of $\mathscr{B}$. This relation is characterized by a rank- 1 factorization of a symmetric $d$-way tensor, constructed from the eigenvalues of the embedding system.

## 2 Notation

The notation used in the paper is standard: $\mathbb{R}$ is the set of real values, $\mathbb{C}$ is the set of complex values, and $\mathbb{N}$ is the set of natural numbers. The set of scalar real-valued signals over $\mathbb{N}$ is denoted by $\mathbb{R}^{\mathbb{N}}$. An autonomous linear time-invariant system $\mathscr{B}$ admits a minimal state space representation

$$
\begin{align*}
& \mathscr{B}=\mathscr{B}(A, c):=\left\{z \in \mathbb{R}^{\mathbb{N}} \mid \text { there is } x,\right. \text { such that } \\
&\left.\sigma x=A x, z=c x, x(1) \in \mathbb{R}^{n}\right\}, \tag{1}
\end{align*}
$$

where $A \in \mathbb{R}^{n \times n}$ and $c \in \mathbb{R}^{1 \times n}$ are parameters of the system and $\sigma$ is the shift operator $(\sigma x)(t)=x(t+1)$. The eigenvalues $\lambda_{1}, \ldots, \lambda_{n}$ of $A$ are invariant of the representation and are, therefore, a property of the system $\mathscr{B}$.

In this paper, we consider a single output Wiener system
and assume that the eigenvalues of its linear subsystem are distinct. In this case, the linear time-invariant subsystem admits a sum-of-damped-exponentials representation

$$
\begin{equation*}
\mathscr{B}=\mathscr{B}(\lambda):=\left\{z \in \mathbb{R}^{\mathbb{N}} \mid z=\sum_{i=1}^{n} \alpha_{i} \exp _{\lambda_{i}}, \alpha \in \mathbb{C}^{n}\right\} \tag{2}
\end{equation*}
$$

where $\exp _{\lambda_{i}}(t):=\lambda_{i}^{t}$ and $\lambda$ is the vector of the system's eigenvalues $\lambda=\left[\begin{array}{lll}\lambda_{1} & \cdots & \lambda_{n}\end{array}\right]^{\top} \in \mathbb{C}^{n}$.

Restricting ourselves to the single output case with distinct eigenvalues simplifies the notation. The results in the paper can be generalized mutatis mutandis to the case of multioutput systems. Dealing with repeated eigenvalues requires a generalization of the sum-of-damped-exponentials representation, which complicates the analysis but does not change our main results.

The static nonlinearity $g$ is a $d$ th order polynomial, represented by a given monomial basis $v$
$y=g(z):=\theta^{\top} v(z), \quad$ where $v(z)=\left[\begin{array}{llll}z^{0} & z^{1} & \cdots & z^{d}\end{array}\right]^{\top} \in \mathbb{R}^{d+1}$.
Putting together (2) and (3), we obtain the autonomous Wiener system

$$
\mathscr{B}_{\mathrm{w}}(\lambda, \theta):=\left\{y \in \mathbb{R}^{\mathbb{N}} \mid(2,3) \text { hold for } \alpha \in \mathbb{C}^{n}\right\}
$$

parameterized by the vector of the coefficients $\theta=$ $\left[\begin{array}{llll}\theta_{0} & \theta_{1} & \cdots & \theta_{d}\end{array}\right]^{\top} \in \mathbb{R}^{d+1}$ of the nonlinear part and the eigenvalues $\lambda$ of the linear part.

The choice of polynomial nonlinearity is motivated by the fact that any smooth function can be approximated arbitrarily well over a finite interval by a polynomial of sufficiently high degree. In this paper, we assume that the degree $d$ of the nonlinearity is a priori given. In practice, however, the question of how to choose it is a relevant (and difficult) one. A general way to estimate this hyper-parameter is to use cross-validation. Computational and sensitivity issues however restrict us to relatively low degree polynomials.

## 3 Main result

Theorem 1 Consider an autonomous Wiener system $\mathscr{B}_{w}(\lambda, \theta)$ with order-n linear subsystem and a degree-d nonlinear subsystem. Assume that the eigenvalues $\lambda$ are distinct and 1 is not an eigenvalue of $\mathscr{B}(\lambda)$. Then, there is an autonomous linear time-invariant system $\mathscr{B}\left(\lambda_{w}\right)$ with eigenvalues $\lambda_{w} \in \mathbb{C}^{n_{w}}$, where

$$
\begin{equation*}
n_{w} \leq \bar{n}_{w}:=\binom{n+d}{d}=\frac{(n+1)(n+2) \cdots(n+d)}{d!} \tag{4}
\end{equation*}
$$

such that

$$
\begin{equation*}
\mathscr{B}_{w}(\lambda, \theta) \subseteq \mathscr{B}\left(\lambda_{w}\right) \tag{5}
\end{equation*}
$$

The eigenvalues $\lambda_{w}$ of the embedding system $\mathscr{B}\left(\lambda_{w}\right)$ are products of $d$ elements of the set $\left\{\lambda_{0}, \lambda_{1} \ldots, \lambda_{n}\right\}$, where $\lambda_{0}:=1$, i.e., there are indices $k_{i, 1}, \ldots, k_{i, d} \in\{0,1, \ldots, n\}$, such that

$$
\begin{equation*}
\lambda_{w, i}=\prod_{j=1}^{d} \lambda_{k_{i, j}}, \quad \text { for } i=1, \ldots, n_{w} . \tag{6}
\end{equation*}
$$

PROOF. By definition

$$
\begin{equation*}
\mathscr{B}\left(\lambda_{\mathrm{w}}\right):=\left\{y \in \mathbb{R}^{\mathbb{N}} \mid y=\sum_{i=1}^{n_{\mathrm{w}}} \beta_{i} \exp _{\lambda_{\mathrm{w}, i}}, \beta \in \mathbb{C}^{n_{\mathrm{w}}}\right\} . \tag{7}
\end{equation*}
$$

In order to prove the relation (5), we compare the output of (7) with the expression for the output of the autonomous Wiener system $\mathscr{B}_{\mathrm{w}}(\lambda, \theta)$.

Consider a general basis element

$$
v_{j}(z(t))=(z(t))^{j}=\left(\sum_{i=1}^{n} \alpha_{i} \lambda_{i}^{t}\right)^{j}
$$

For $j=0$ and $1, v_{0}=1$ and $v_{1}=y$ are of the form of sum-of-damped-exponentials with $n_{0}=1$ and $n_{1}=n$ exponents, respectively. For $j>1, v_{j}$ is also of the form of a sum-of-damped-exponentials with exponents that are products of $j$ elements of the set $\lambda$, i.e.,

$$
v_{j}(z(t))=\sum_{i=1}^{n_{j}} \gamma_{i} \mu_{i, j}^{t}, \quad \text { where } \mu_{i, j}^{t}=\prod_{\ell=1}^{j} \lambda_{k_{i, j, \ell}}
$$

for some indices $k_{i, j, \ell} \in\{1, \ldots, n\}$. The number of terms $n_{j}$ is equal to the number of combinations with repetitions of $j$ elements out of the $n$ elements of $\lambda$. Therefore,

$$
n_{j}=\binom{n+j-1}{j}=\frac{(n+j-1) \cdots n}{j!} .
$$

Consider now the output

$$
y(t)=g(z(t))=\theta^{\top} v(z(t)) .
$$

It is also of the form of a sum-of-damped-exponentials

$$
\begin{equation*}
y(t)=\sum_{i=1}^{n_{\mathrm{w}}} \zeta_{i} \lambda_{\mathrm{w}, i}^{t}, \quad \text { where }\left\{\lambda_{\mathrm{w}, 1}, \ldots, \lambda_{\mathrm{w}, n_{\mathrm{w}}}\right\}=\bigcup_{i=0}^{j} \bigcup_{j=0}^{d} \mu_{i, j} . \tag{8}
\end{equation*}
$$

The elements of $\lambda_{\mathrm{w}}$ are products of $d$ elements of the set $\left\{1, \lambda_{1}, \ldots, \lambda_{n}\right\}$. The number of such products is

$$
\bar{n}_{\mathrm{w}}=\sum_{j=0}^{d} n_{j}=\binom{n+d}{d}=\frac{(n+1)(n+2) \cdots(n+d)}{d!} .
$$

The number of distinct elements $n_{\mathrm{w}}$, which is the order of $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$, is therefore upper bounded by $\bar{n}_{\mathrm{w}}$.

We've shown that both the output (7) of the autonomous Wiener system $\mathscr{B}_{\mathrm{w}}(\lambda, \theta)$ and the output (8) of the linear system $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$ are of the form of sum-of-damped-exponentials with the same exponents. The coefficients $\beta_{i}$ in (7) however are restricted only by the condition of being in complex conjugate pairs (since the signal is real), while the coefficients $\zeta_{i}$ in (8) range over an $n$-dimensional manifold of $\mathbb{C}^{n_{\mathrm{w}}}$. This proves (5).

Note 1 The generalization of Theorem 1 to multi-output systems can be done by considering each output separately and including polynomial factors in (7). Another approach is to use a state space representation, in which case the sum-ofexponentials expression in (7) becomes $y(t)=C A^{t} x(0)$.

Next, we give an alternative characterization of the relation (6) between the eigenvalues of $\lambda$ and $\lambda_{\mathrm{w}}$, using the notation " $\circ$ " for the vector outer product.

Corollary 2 (Link between $\lambda_{\mathrm{w}}$ and $\lambda$ ) The symmetric, rank-1, $d$-way tensor $T:=\underbrace{\lambda \circ \lambda \circ \cdots \circ \lambda}_{d \text { times }}$, has as unique elements $\lambda_{w, 1}, \ldots, \lambda_{w, n_{w}}$.

## 4 Application of the result in system identification

The problem considered in this section is: Given a monomial basis $v$, a finite trajectory $y_{\mathrm{d}}=\left(y_{\mathrm{d}}(1), \ldots, y_{\mathrm{d}}(T)\right)$ of an autonomous Wiener system $\mathscr{B}_{\mathrm{w}}(\lambda, \theta)$, and the order $n$ of its linear part, find parameters $\widehat{\lambda}, \widehat{\theta}$, such that

$$
\mathscr{B}_{\mathrm{w}}(\lambda, \theta)=\mathscr{B}_{\mathrm{w}}(\widehat{\lambda}, \widehat{\theta})
$$

Theorem 1 suggests the following solution method:
(1) identify the embedding system $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$ from $y_{\mathrm{d}}$,
(2) compute the linear subsystem $\mathscr{B}(\lambda)$ from $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$, and
(3) compute the nonlinear subsystem from $\mathscr{B}\left(\lambda_{\mathrm{w}}\right), \mathscr{B}(\lambda)$.

Assuming that the given trajectory $y_{\mathrm{d}}$ is persistently exciting of order $n_{\mathrm{w}}$, the embedding system $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$ is identifiable from $y_{\mathrm{d}}$ [19]. The remaining problems, resolved in steps 2 and 3 , are to find from the identified system $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$, the linear and nonlinear subsystems of the autonomous Wiener system. Note that due to exchange of gain, between $\mathscr{B}$ and $g$, the linear and nonlinear subsystems are not identifiable from the data alone. As shown next, however, the eigenvalues of $\mathscr{B}$ can be determined uniquely and $g$ can be determined up to a scaling factor.

Note 2 Classical identifiability conditions, see, e.g., [10,9,19], require persistency of excitation of the input sequence. They are not applicable to autonomous systems. In
this paper, the persistency of excitation assumption is on the output sequence. Similar assumption is used in dynamic measurement [16,11,12], where the data is also a response of an autonomous system.

### 4.1 Identification of $\mathscr{B}\left(\lambda_{w}\right)$ from the given output data

The identification of an autonomous linear time-invariant system $\mathscr{B}_{\mathrm{w}}(\lambda, \theta)$ from the finite trajectory $y_{\mathrm{d}} \in \mathscr{B}_{\mathrm{w}}(\lambda, \theta)$ is a classical problem, see for example, $[7,6]$ and $[13$, Section 5.1.3]. One possible solution [6] is to form the Hankel matrix
$\mathscr{H}_{n_{\mathrm{w}}+1}\left(y_{\mathrm{d}}\right):=\left[\begin{array}{cccc}y_{\mathrm{d}}(1) & y_{\mathrm{d}}(2) & \cdots & y_{\mathrm{d}}\left(T-n_{\mathrm{w}}\right) \\ y_{\mathrm{d}}(2) & y_{\mathrm{d}}(3) & \cdots & y_{\mathrm{d}}\left(T-n_{\mathrm{w}}+1\right) \\ \vdots & \vdots & & \vdots \\ y_{\mathrm{d}}\left(n_{\mathrm{w}}+1\right) & y_{\mathrm{d}}\left(n_{\mathrm{w}}+2\right) & \cdots & y_{\mathrm{d}}(T)\end{array}\right]$
and compute its left kernel (which can be shown to be one dimensional)

$$
p \mathscr{H}_{n_{\mathrm{w}}+1}\left(y_{\mathrm{d}}\right)=0 .
$$

The roots of the polynomial $p(s)=p_{0}+p_{1} s+\cdots+p_{n_{\mathrm{w}}} s^{n_{\mathrm{w}}}$ are the eigenvalues of the embedding system. Another solution (called Kung's method [7]) is based on realization theory: 1) compute the rank revealing factorization

$$
\mathscr{H}_{n_{\mathrm{w}}+1}\left(y_{\mathrm{d}}\right)=\mathscr{O} \mathscr{C}, \quad \text { with } \mathscr{O} \in \mathbb{R}^{L \times n_{\mathrm{w}}} \text { and } \mathscr{C}^{n_{\mathrm{w}} \times(T-L)}
$$

of the Hankel matrix $\mathscr{H}_{L}\left(y_{\mathrm{d}}\right)$, where $L$ is a design parameter, satisfying the constraints $n_{\mathrm{w}}+1 \leq L \leq T-n_{\mathrm{w}}$ and 2) solve the system of linear equations $\overline{\mathscr{O}} \widehat{A}=\underline{\mathscr{O}}$, for $\widehat{A}$, where $\overline{\mathscr{O}}$ is the matrix $\mathscr{O}$ with the first row removed and $\underline{O}$ is the matrix $\mathscr{O}$ with the last row removed. The eigenvalues of $\widehat{A}$ are the eigenvalues of the embedding system.

The minimal number of samples needed for the identification of the system $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$ is $T_{\text {min }}=2 n_{\mathrm{w}}+1$. The identification data, however, can be collected from $n_{\mathrm{w}}$ experiments with $n_{\mathrm{w}}+1$ samples instead of a single experiment with $T_{\text {min }}$ samples. Let $y_{\mathrm{d}}^{1}, \ldots, y_{\mathrm{d}}^{n_{\mathrm{w}}}$ be the data of the multiple experiments of length $n_{\mathrm{w}}+1$. Then, the identification procedure is modified by replacing the Hankel matrix $\mathscr{H}_{n_{\mathrm{w}}+1}\left(y_{\mathrm{d}}\right)$ by the matrix $\left[y_{\mathrm{d}}{ }^{1} \cdots y_{\mathrm{d}}{ }^{n_{\mathrm{w}}}\right]$ of the stacked next to each other responses. More generally, using date from multiple experiments of length $T_{1}, \ldots, T_{n_{\mathrm{w}}}>n_{\mathrm{w}}$, the identification method is based on the computation of the left kernel or rank revealing factorization of the mosaic Hankel matrix [14]

$$
\mathscr{H}_{n_{\mathrm{w}}+1}\left(y_{\mathrm{d}}^{1}, \ldots, y_{\mathrm{d}}^{n_{\mathrm{w}}}\right):=\left[\mathscr{H}_{n_{\mathrm{w}}+1}\left(y_{\mathrm{d}}^{1}\right) \cdots \mathscr{H}_{n_{\mathrm{w}}+1}\left(y_{\mathrm{d}}^{n_{\mathrm{w}}}\right)\right] .
$$

### 4.2 Computation of the subsystem $\mathscr{B}(\lambda)$ from $\mathscr{B}\left(\lambda_{w}\right)$

After finding $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$, the next step is the computation of the linear subsystem. We are interested in the transition from
$\lambda_{\mathrm{w}}$ to $\lambda$, i.e., extracting the linear subsystem $\mathscr{B}(\lambda)$ from $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$. Using Corollary 2 , we can find $\lambda$ by computing a rank-1 factorization of a symmetric tensor $T$ constructed from $\lambda_{\mathrm{w}}$. Checking whether $T$ has rank equal to one can be done by checking the rank of the $d$ unfoldings of the tensor: $T$ is rank-1 if and only if all unfoldings of $T$ are rank-1 [3].

Another characterization of (6) that leads to a more efficient method is given in terms of the "frequencies" $\omega_{i}:=\angle \lambda_{i}$ and $\omega_{\mathrm{w}, i}:=\angle \lambda_{\mathrm{w}, i}$ of $\mathscr{B}(\lambda)$ and $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$, respectively. From (6), we have the following linear relation among $\omega_{\mathrm{w}, i}$ 's and $\omega_{i}$ 's

$$
\omega_{\mathrm{w}, i}=\sum_{j=1}^{d} \omega_{k_{i, j}} \quad(\bmod 2 \pi)
$$

Therefore, there is an $n_{\mathrm{w}} \times n$ matrix $K$, such that

$$
\begin{equation*}
\omega_{\mathrm{w}}=K \omega \quad(\bmod 2 \pi) \tag{9}
\end{equation*}
$$

Relation (9) shows that the problem of extracting $\mathscr{B}(\boldsymbol{\lambda})$ from $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$ can be solved by computing the frequencies of $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$ and solving a system of linear equations. The ordering of the $\omega_{\mathrm{w}}$ 's however is unknown, so that all permutations of the $\omega_{\mathrm{w}}$ 's should be tested for existence of an exact solution. (The order of the to-be-found frequencies $\omega$ is not important.) This method requires the same number of subproblems to-be-solved as in the procedure using Corollary 2. The subproblem (9) however is a linear system, which is simpler and faster to solve than the rank-1 factorization of a symmetric tensor.

### 4.3 Computation of the subsystem $g$ from $\mathscr{B}\left(\lambda_{w}\right), \mathscr{B}(\lambda)$

Finally computing the nonlinear function $g$ requires a simultaneous rank-1 factorization of $d$ tensors. General theory and methods, called structured data fusion, for solving simultaneous tensor factorization problems is developed in [17]. In general, the structured data fusion problem has no analytical solution and requires iterative solution methods. Applied to the autonomous Wiener system identification problem, however, when $g$ contains a first and/or second order terms the structured data fusion problem has a trivial solution. When $g$ has a first order term, the coefficients $\theta$ can be obtained directly from the coefficients $\gamma$ in (8) without extra computations. When $g$ has a second order term, the coefficients $\theta$ can be obtained from the coefficients $\gamma$ in (8) by a Cholesky factorization of a symmetric matrix constructed from the $\gamma$ 's.

Note 3 The method presented in this section is for exact identification. In order to show its performance in the case of measurement noise, in Section 5 we show results of a Monte-Carlo simulation. Another issue for practical application of the method is robustness to modeling errors. Although this is a difficult problem to analyze theoretically, there are empirical tests that can indicate the presence of modeling errors in an identified model, e.g., assuming that
the measurement noise is white, colored residuals indicate the presence of modeling errors, see [8, Ch. 16, page 512].

## 5 Numerical example

The autonomous Wiener system $\mathscr{B}_{\mathrm{w}}(\lambda, \theta)$ used in the simulation example consists of a second order linear subsystem with eigenvalues $\lambda_{1,2}=-0.5 \pm 0.7 \mathbf{i}$ and a (dead zone) nonlinear subsystem defined by the third order polynomial

$$
g(z)=\theta_{0}+\theta_{1} z+\theta_{2} z^{2}+\theta_{3} z^{3}
$$

with coefficients $\theta=\left[\begin{array}{llll}1 & 1 & 1 & 1\end{array}\right]^{\top}$.
According to Theorem $1, \mathscr{B}_{\mathrm{w}}(\lambda, \theta)$ is included in a linear time-invariant system of order

$$
\begin{equation*}
n_{\mathrm{w}} \leq\binom{ n+d}{d}=\binom{5}{3}=10 \tag{10}
\end{equation*}
$$

In order to verify this property empirically, we generate a $T=25$ samples long trajectory $y_{\mathrm{d}}$ of the system $\mathscr{B}_{\mathrm{w}}(\lambda, \theta)$ due to a random initial condition and check the rank of the square Hankel matrix $\mathscr{H}_{13}\left(y_{\mathrm{d}}\right)$. The fact that $\operatorname{rank}\left(\mathscr{H}_{13}\left(y_{\mathrm{d}}\right)\right)=10$ confirms (10).

Next, we verify (6), namely the statement that the eigenvalues $\lambda_{\mathrm{w}}$ of the embedding system are products of up to $d$ eigenvalues $\lambda$ of the linear subsystem. First, using an exact identification method, e.g., Kung's method described in Section 4, we obtain the linear time-invariant system $\mathscr{B}\left(\lambda_{\mathrm{w}}\right)$ that contains $\mathscr{B}_{\mathrm{w}}(\lambda, \theta)$. Then, we form the set

$$
\begin{equation*}
\left\{\lambda_{k_{0}} \lambda_{k_{1}} \cdots \lambda_{k_{d}} \mid \lambda_{0}:=1 \text { and } k_{0}, k_{1}, \ldots, k_{d} \in\{0,1 \ldots, n\}\right\} \tag{11}
\end{equation*}
$$

of all products of up to $d$ eigenvalues of $\mathscr{B}$. Finally, we compare the identified eigenvalues $\lambda_{\mathrm{w}}$ and the theoretically predicted ones (11). For the simulation example, described above, we confirm that they coincide, see Figure 1.


Fig. 1. The eigenvalues $\lambda_{\mathrm{w}}$ (plotted as red $\times$ 's) of the embedding system coincide with the products (plotted as blue o's) of up to $d$ eigenvalues of the linear subsystem $\mathscr{B}(\lambda)$. + - eigenvalues $\lambda$, dotted line - unit circle.
For higher values of $n$ and $d$ it is possible to obtain after identification from data only a subset of the eigenvalues of $\mathscr{B}_{\mathrm{w}}$. This is due to ill-conditioning of the identification problem and the finite precision arithmetic used in the numerical computations. The problem can be partially resolved using
data of multiple experiments generated by properly selected initial conditions (design of the experiments). However, the ill-conditioning of the system identification problem remains an important practical issue that will be addressed elsewhere.

In order to study the robustness of $\hat{\lambda}_{w}$ to measurement noise, next, we add to the exact output $y$ white Gaussian noise with standard deviation $s=0.25$. Figure 2 shows the exact eigenvalues of the embedding system ( $\circ$ ) and the computed ones $(\times)$ by Kung's method for 25 different noise realizations. The scatter plots of the estimated eigenvalues indicates the sensitivity of $\hat{\lambda}_{\mathrm{w}}$ to the noise. Eigenvalues close to the unit circle seem to more difficult to estimate as their estimates have large variances.

Particularly difficult to estimate is the eigenvalue at 1 , which comes from the constant term in the static nonlinearity.


Fig. 2. o-true $\lambda_{\mathrm{w}}, \times-25$ estimates $\hat{\lambda}_{\mathrm{w}}$, obtained with Kung's method. Eigenvalues that are close to the unit circle are more sensitive to additive noisy in the output data. The poorly estimated poles are the ones for $\lambda_{\mathrm{w}, i}=1$.

We define the relative parameter estimation error

$$
e=\|\bar{\theta}-\widehat{\theta}\| /\|\bar{\theta}\|
$$

where $\bar{\theta}=\left[\begin{array}{llll}1 & 1 & 1 & 1\end{array}\right]^{\top}$ is the true system's parameter vector and $\widehat{\theta}$ is the normalized estimated parameter vector by Kung's method. (In order to avoid the nonuniqueness of $\widehat{\theta}$, we impose the normalization $\widehat{\theta}_{0}=1$.)

Figure 3 shows the averaged over 500 Monte-Carlo repetitions parameter estimation error as a function of the noise-to-signal ratio $(=1 / \mathrm{SNR})$. The result shows a typical for subspace methods gradual increase of the error. After a threshold noise level, the increase of the error is faster. The subspace method "works well" up this threshold value. In the numerical example shown in Figure 3 the threshold is around $\mathrm{SNR}=77 \mathrm{~dB}$, which is achieved only in high-quality measurement environments. The result suggests that the subspace approach is sensitive to disturbances making it impractical for real-life applications.

The high sensitivity of the estimated model to noise is due to the fact that the system is of high order $\left(n_{\mathrm{w}}=10\right)$ relative to the length of the data $(T=25)$. Since the data is collected from a transient response, archiving better results by measuring longer is not an option. This leads to an ill-conditioning of the problem, an issue that is independent of a particular method being used for solving the problem. As explained in Section 4.1, one way to avoid the ill-conditioning issue is to measure data from multiple experiments. Also, the sensitivity of the subspace method can be reduced by using extra prior knowledge about the spectrum of the true system, in particular the fixed eigenvalue at 1 .


Fig. 3. The method delivers exact result for noise-free data. The estimation error increases as a function of the noise-to-signal ratio. After a threshold SNR, the increase is faster. The method "works well" up this threshold value. Without using extra prior knowledge about the spectrum the threshold is around 77 dB .

## 6 Conclusion

We showed that the behavior of an autonomous Wiener system with polynomial nonlinearity is included in the behavior of a finite-dimensional linear time-invariant system. The order of the embedding linear system depends combinatorially on the order of the linear subsystem and the degree of the static nonlinearity. The relation between the eigenvalues of the embedding system and the linear subsystem is given by a rank- 1 factorization of a symmetric tensor: the unique elements of the tensor are the eigenvalues of the embedding system and the factors contain the eigenvalues of the linear subsystem. The result suggests an autonomous Wiener system identification procedure that is based on linear timeinvariant system identification followed by solution of a system of linear equations.

Challenges that need to be addressed in order to make this procedure practically useful are ill-conditioning of the linear identification step and combinatorial number of systems that have to be solved for the computation of the eigenvalues of the linear subsystem from the eigenvalues of the identified system. Future work includes modification and extension of the method to approximate identification of autonomous Wiener systems, statistical analysis of the resulting method in the output error setup, and generalization of the results to Wiener systems with inputs as well as initial conditions.

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