On backward shift algorithm for estimating poles of systems \star

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Abstract

In this paper, we present an algorithm of estimating poles of linear time-invariant systems by using the backward shift operator. It is proved that poles of rational functions, including zeros and multiplicities, are solutions of an algebra equation which can be obtained by taking backward shift operator to normalized reproductive kernels in the unit disc case. The algorithm is accordingly developed for frequency-domain identification. The robustness of this algorithm is proved. Some illustrative examples are presented to show the efficiency for systems with distinguished and multiple poles cases.

 $Key \ words:$ Backward shift operator, poles estimation, rational orthogonal basis, linear time-invariant systems, identification methods

1 Introduction

System identification is to build mathematical models which fit the measured data from discrete or continuous systems. An amount of methods have been developed for this problem, such as [5,6,26-28]. A classical guidebook for one getting to know this topic is [9]. For system identification of linear time-invariant (LTI) systems, the priori-knowledge of poles is important, especially for the methods that adopt rational orthogonal bases such as in[15–17,25]. In these methods, the estimated poles are used to construct rational orthogonal basis functions. A collection of these excellent results is [7].

In unit disc case, the general setting of a rational orthogonal basis is

$$\mathcal{B}_{k}(z) = \mathcal{B}_{\{a_{1},.a_{k}\}}(z) \triangleq \frac{\sqrt{1 - |a_{k}|^{2}}}{1 - \overline{a}_{k}z} \prod_{l=1}^{k-1} \frac{z - a_{l}}{1 - \overline{a}_{l}z}, \quad (1)$$

where a_k s (k = 1, ...) are in the unit disc, $(\overline{a} \text{ means conjugation of } a)$. Many researchers work on choosing optimal n-poles $\{a_k\}_{k=1}^n$ in order to define the best rational orthogonal bases for a system. T. Oliveria e Silva derived

the optimal pole conditions for Laguerre, Kautz and general orthogonal basis function models in [18–20], respectively. In [10–12], adaptive selection of poles is studied. Other attempts to estimate optimal pole positions of a Laguerre model are given in [23,3]. Generally speaking, the pole estimation of an LTI system, in practice, is not easy.

For a discrete LTI system which is causal and stable, let $\{x_k\}, \{y_k\}$ be the input and output signals, respectively. There is a relation between $\{x_k\}$ and $\{y_k\}$ as

$$y_k = \{x_k\} * \{h_k\} = \sum_{l=0}^{+\infty} h_l x_{k-l}, \qquad (2)$$

where $\{h_k\}$ is the impulse response. With an operator q, qx(k) = x(k+1), is drawn into, (2) can be represented as

$$y_k = \sum_{l=0}^{+\infty} h_l x_{k-l} = (\sum_{l=0}^{+\infty} h_l q^{-l}) x_k.$$
 (3)

The related function

$$G(z) = \sum_{l=0}^{+\infty} h_l z^{-l} \tag{4}$$

is the transfer function of the system. The values of the transfer function for z on the unit circle are called frequency responses. Under the stability and casuality as-

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sumption, G(z) is a proper rational function with real coefficients. To identify G(z) with general orthogonal bases, estimating poles for the basis functions plays a significant role.

As is well-known, the basis (1) can be obtained by the normalized reproducing kernels $e_a(z)$ with Gram-Schmidt process, where the normalized reproducing kernel at a, is given by

$$e_a(z) = \frac{\sqrt{1 - |a|^2}}{1 - \overline{a}z}$$

For these kernels, there is a very good property when taking backward shift operator on them. Based on this property, we can estimate poles of e_a 's instead of the orthogonal cases. In this paper, we are to locate poles of an LTI system based on a set of frequency domain measurements by using backward shift operator, which results in an algorithm, we call it *backward shift algorithm*.

This paper is arranged as follows. In section 2, we study each case of taking backward shift operator to the rational functions. After that, we introduce the backward shift algorithm in detail in section 3. Examples are given in section 4 to illustrate the proposed idea. Some conclusions are drawn in the last section.

2 Backward shift on rational functions

2.1 Backward shift operator

The backward shift operator, denoted by \mathbf{S} ,

$$\mathbf{S}(f)(z) = \frac{f(z) - f(0)}{z},\tag{5}$$

is the Banach space adjoint of the forward shift operator $\mathbf{F}(f)(z) = zf(z)$ in the Hardy-2 space in the unit disc, viz.

$$\langle \mathbf{S}(f), g \rangle = \langle f, \mathbf{F}(g) \rangle, \ f, g \in H_2.$$
 (6)

It is an important and interesting operator. Comprehensive studies in the operator and related topics can be found, for instance, in [1,14,4]. It is well known that a collection of countably many reproducing kernels of the Hardy space H_2 , viz., conjugates of shifted Cauchy kernels, generates a backward shift invariant subspace.

For $0 \neq a \in \mathbb{D}$, the unit disc, we notice the kernel $e_a(z) = \frac{1}{1-\overline{a}z}$ (For convenience we will call \overline{a} a pole of it, although we know precisely it is $\frac{1}{\overline{a}}$) is an eigenvector of **S**, viz.,

$$\mathbf{S}(e_a)(z) = \frac{e_a(z) - e_a(0)}{z}$$
$$= \frac{\overline{a}}{1 - \overline{a}z}.$$

Therefore,

$$\begin{split} \mathbf{S}^2(e_a)(z) &= \mathbf{S}(\mathbf{S}(e_a))(z) \\ &= \frac{\overline{a}^2}{1 - \overline{a}z}, \end{split}$$

and, in general,

$$\mathbf{S}^{n}(e_{a})(z) = \frac{\overline{a}^{n}}{1 - \overline{a}z}.$$
(7)

An n-tuple $(a_1, ..., a_n)$ in the unit disc correspond to one of the following two n-tuples of partial fractions, being determined on whether some a_k 's are zero. Denote by $b_1, ..., b_m$ all the distinguished ones among $a_1, ..., a_n$.

Case 1.

If none of the distinguished b_k 's is zero, then it corresponds to

$$\frac{1}{1-\overline{b}_1 z},...,\frac{1}{(1-\overline{b}_1 z)^{l_1}},...,\frac{1}{1-\overline{b}_m z},...,\frac{1}{(1-\overline{b}_m z)^{l_m}},$$

where $l_1, ..., l_m$ are multiples of $b_1, ..., b_m$, respectively and $l_1 + ... + l_m = n$.

A rational function p/q, where p and q are co-prime polynomials, is a non-degenerate linear combination of the above linearly independent set of functions if and only if the degree of q is equal to n, and the degree of p is less than n.

Case 2.

If one of the distinguished b_k 's is zero, say, $b_1 = 0$, with multiplicity l_1 , then it corresponds to

$$1, ..., z^{l_1}, \frac{1}{1 - \overline{b}_2 z}, ..., \frac{1}{(1 - \overline{b}_2 z)^{l_2}}, ..., \frac{1}{(1 - \overline{b}_m z)^{l_m}},$$

where $l_1 + ... + l_m = n$.

A rational function p/q, where p and q are co-prime polynomials, is a non-degenerate linear combination of the above linearly independent set of functions if and only if the degree of q is equal to $n - l_1$, and the degree of p is less than n.

These cases will be studied in detail in the following three subsections.

2.2 The distinguished non-zero poles case

In this subsection we treat the case where all b_k , k = 1, ..., n, are different with each other, that is, each mul-

tiplicity is 1. Assume that f is of the form

$$f(z) = \sum_{k=1}^{n} \frac{\lambda_k}{1 - \overline{b}_k z},\tag{8}$$

where λ_k s are non-zero, b_k s are non-zero and distinguished from each other. Applying, consecutively, the backward shift operator **S** to f(z) *n* times, we have

$$\begin{cases} \mathbf{S}(f)(z) = \frac{\lambda_1 \overline{b}_1}{1 - \overline{b}_1 z} + \frac{\lambda_2 \overline{b}_2}{1 - \overline{b}_2 z} + \ldots + \frac{\lambda_n \overline{b}_n}{1 - \overline{b}_n z}, \\ \mathbf{S}^2(f)(z) = \frac{\lambda_1 \overline{b}_1^2}{1 - \overline{b}_1 z} + \frac{\lambda_2 \overline{b}_2^2}{1 - \overline{b}_2 z} + \ldots + \frac{\lambda_n \overline{b}_n^2}{1 - \overline{b}_n z}, \\ \vdots \\ \mathbf{S}^n(f)(z) = \frac{\lambda_1 \overline{b}_1^n}{1 - \overline{b}_1 z} + \frac{\lambda_2 \overline{b}_2^n}{1 - \overline{b}_2 z} + \ldots + \frac{\lambda_n \overline{b}_n^n}{1 - \overline{b}_n z}. \end{cases}$$

Since the b_k s are distinguished, $\{\frac{1}{1-\overline{b}_{kz}}\}_{k=1}^n$ is a linearly independent collection. There exists a unique non-zero sequence $\{\mu_k\}_{k=0}^n$ such that

$$\mu_0 f(z) + \mu_1 \mathbf{S}(f)(z) + \dots + \mu_n \mathbf{S}^n(f)(z) = 0.$$
 (9)

Precisely,

$$\begin{aligned} f & 0 = (\mu_0 + \mu_1 \overline{b}_1 + \dots + \mu_{n-1} \overline{b}_1^{n-1} + \mu_n \overline{b}_1^n) \frac{\lambda_1}{1 - \overline{b}_1 z} \\ & + (\mu_0 + \mu_1 \overline{b}_2 + \dots + \mu_{n-1} \overline{b}_2^{n-1} + \mu_n \overline{b}_2^n) \frac{\lambda_2}{1 - \overline{b}_2 z} \\ & \vdots \\ & + (\mu_0 + \mu_1 \overline{b}_n + \dots + \mu_{n-1} \overline{b}_n^{n-1} + \mu_n \overline{b}_n^n) \frac{\lambda_n}{1 - \overline{b}_n z}. \end{aligned}$$

The linear independence of $\{\frac{1}{1-\bar{b}_kz}\}_{k=1}^n$ implies

$$\begin{cases} \mu_0 + \mu_1 \overline{b}_1 + \mu_2 \overline{b}_1^2 + \dots + \mu_{n-1} \overline{b}_1^{n-1} + \mu_n \overline{b}_1^n = 0, \\ \mu_0 + \mu_1 \overline{b}_2 + \mu_2 \overline{b}_2^2 + \dots + \mu_{n-1} \overline{b}_2^{n-1} + \mu_n \overline{b}_2^n = 0, \\ \vdots \\ \mu_0 + \mu_1 \overline{b}_n + \mu_2 \overline{b}_n^2 + \dots + \mu_{n-1} \overline{b}_n^{n-1} + \mu_n \overline{b}_n^n = 0, \end{cases}$$

Which means b_k , k = 1, ..., n, are the solutions of an algebraic equation

$$\mu_n x^n + \mu_{n-1} x^{n-1} + \mu_{n-2} x^{n-2} + \dots + \mu_1 x + \mu_0 = 0.$$
 (10)

The equation (10) has n different roots, then $\mu_n \neq 0$, we may assume $\mu_n = 1$.

2.3 The Multiple But Non-Zero Poles Case

In this subsection, we deal with rational function of order n with non-zero multiple poles. Without loss of generality, we assume f(z) to be of the form

$$f(z) = \sum_{j=1}^{m} \frac{\lambda_j}{(1 - \bar{b}_1 z)^j} + \sum_{k=m+1}^{n} \frac{\lambda_k}{1 - \bar{b}_k z},$$
 (11)

where λ_k , for $k \geq m$, are all nonzero. First of all, by applying the backward shift operator to $\frac{1}{(1-\overline{b}_1 z)^m}$, we have

$$\mathbf{S}\left(\frac{1}{(1-\bar{b}_{1}z)^{m}}\right) = \bar{b}_{1}\left(\frac{1}{(1-\bar{b}_{1}z)^{m}} + \frac{1}{(1-\bar{b}_{1}z)^{m-1}} + \dots + \frac{1}{(1-\bar{b}_{1}z)}\right).$$
(12)

Set
$$M(z) = \sum_{j=1}^{m} \frac{1}{(1-\overline{b}_1 z)^j}$$
, then

$$\begin{aligned} \mathbf{S}(M)(z) \\ &= \bar{b}_1 \frac{1}{(1 - \bar{b}_1 z)} + \bar{b}_1 \left(\frac{1}{(1 - \bar{b}_1 z)^2} + \frac{1}{(1 - \bar{b}_1 z)} \right) \\ &+ \bar{b}_1 \left(\frac{1}{(1 - \bar{b}_1 z)^m} + \dots \frac{1}{(1 - \bar{b}_1 z)} \right) \\ &= \bar{b}_1 \left(\frac{1}{(1 - \bar{b}_1 z)^m} + \frac{2}{(1 - \bar{b}_1 z)^{m-1}} + \dots + \frac{m}{(1 - \bar{b}_1 z)} \right). \end{aligned}$$
(13)

Recursively,

$$\mathbf{S}^{k}(M)(z) = \overline{b}_{1}^{k} \left(\frac{c_{m}^{k}}{(1 - \overline{b}_{1}z)^{m}} + \frac{c_{m-1}^{k}}{(1 - \overline{b}_{1}z)^{m-1}} + \dots + \frac{c_{1}^{k}}{(1 - \overline{b}_{1}z)} \right),$$
(14)

where c_j^k , j = 1, ..., m; k = 0, 1, 2, ..., are the coefficients of the order-jth term by taking kth backward shift operator. We can see the coefficients also satisfy the recursive rule,

$$c_j^{k+1} = \sum_{l=j}^{m} c_l^k, (j = 1, 2, ..., m).$$
(15)

Obviously $c_m^0 = c_{m-1}^0 = \dots = c_1^0 = 1$, from (12) we know by taking backward shift operator to $\frac{1}{(1-\overline{b}_1z)^m}$, it will generate m components in which the coefficients of the m-th order is 1, then there is $c_m^1 = c_m^2 \dots = c_m^k = c_m^{k+1} = \dots = 1$. From (13), due to linear independence of $\{\frac{1}{(1-\overline{b}_1z)^j}, j = 1, \dots, m; \frac{1}{1-\overline{b}_{m+1}z}, \dots, \frac{1}{1-\overline{b}_nz}\}$, there exists a unique sequence $\{\mu_k\}_{k=0}^n$, such that the same relation

as equation (9) holds,

$$\mu_0 f(z) + \mu_1 \mathbf{S}(f)(z) \dots + \mu_n \mathbf{S}^n(f)(z) = 0.$$
 (16)

That results to

$$\begin{cases} \mu_{0} + \mu_{1}\overline{b}_{1} + \mu_{2}\overline{b}_{1}^{2} + \dots + \mu_{n}\overline{b}_{1}^{n} = 0, \\ \mu_{0} + \mu_{1}\overline{b}_{m+1}\dots + \mu_{n}\overline{b}_{m+1}^{n} = 0, \\ \vdots \\ \mu_{0} + \mu_{1}\overline{b}_{n} + \mu_{2}\overline{b}_{n}^{2} + \dots + \overline{b}_{n}^{n} = 0, \end{cases}$$
(17)

and equations about \overline{b}_1 ,

$$\begin{aligned} c_m^0 \mu_0 + c_m^1 \mu_1 \bar{b}_1 + \dots c_m^n \mu_n \bar{b}_1^n &= 0, \\ c_{m-1}^0 \mu_0 + c_{m-1}^1 \mu_1 \bar{b}_1 + \dots + c_{m-1}^n \mu_n \bar{b}_1^n &= 0, \\ &\vdots \\ c_2^0 \mu_0 + c_2^1 \mu_1 \bar{b}_1 + \dots c_2^n \mu_n \bar{b}_1^n &= 0, \\ c_1^0 \mu_0 + c_1^1 \mu_1 \bar{b}_1 + \dots c_1^n \mu_n \bar{b}_1^n &= 0. \end{aligned}$$

There are m equations of \bar{b}_1 . By first eliminating μ_0 we have m-1 equations given by,

$$\begin{cases} c_{m-1}^{0}\mu_{1} + c_{m-1}^{1}\mu_{2}\overline{b}_{1} + \dots c_{m-1}^{n-1}\mu_{n}\overline{b}_{1}^{n-1} = 0, \\ c_{m-2}^{0}\mu_{1} + c_{m-2}^{1}\mu_{2}\overline{b}_{1} + \dots c_{m-2}^{n-1}\mu_{n}\overline{b}_{1}^{n-1} = 0, \\ \vdots \\ c_{2}^{0}\mu_{1} + c_{2}^{1}\mu_{2}\overline{b}_{1} + \dots c_{2}^{n-1}\mu_{n}\overline{b}_{1}^{n-1} = 0, \\ c_{1}^{0}\mu_{1} + c_{1}^{1}\mu_{2}\overline{b}_{1} + \dots c_{1}^{n-1}\mu_{n}\overline{b}_{1}^{n-1} = 0. \end{cases}$$

Repeating this process (m-1) times, it comes

$$c_1^0 \mu_{m-1} + c_1^1 \mu_m \bar{b}_1 + \dots c_1^{n-m+1} \mu_n \bar{b}_1^{n-m+1} = 0.$$
(18)

The following is a result about the coefficients c_1^k .

Lemma 1 For any positive integer m, let $c_m^0 = c_{m-1}^0 = \dots = c_1^0 = \dots = c_m^k = c_m^{k+1} = \dots = 1$, and c_j^k be given by the recursive formula (15), then we have

$$c_1^k = \frac{(m+k-1)!}{(m-1)!k!}.$$
(19)

Proof. We use mathematical induction. First of all, $c_1^0 = \frac{(m-1)!}{(m-1)!0!} = 1$ and $c_1^1 = \frac{(m+1-1)!}{(m-1)!1!} = m$. Second, assume that the result holds for c_1^l , $1 < l \leq k$. For l = k + 1, observing that c_l^k (l = 1, ..., m) are the numbers of

combinations, we have

$$c_{1}^{k} = C_{m+k-1}^{k}, \\ c_{2}^{k} = C_{m+k-2}^{k}, \\ \vdots \\ c_{m}^{k} = C_{k}^{k}.$$

Then there holds

$$c_1^{k+1} = \sum_{l=1}^m c_l^k = \sum_{l=1}^m C_{m+k-l}^k.$$

Using the property of combination $C_m^n + C_m^{n+1} = C_{m+1}^{n+1}$ (for any positive integers m, n), we obtain

$$c_1^{k+1} = C_{m+k}^{k+1} = \frac{(m+k)!}{(m-1)!(k+1)!}.$$
 (20)

The proof is complete.

The first equation in (17) implies \bar{b}_1 is a solution of an equation like (10). Equation (18) is the (m-1)th derivative of the first equation in (17) by lemma 1, then \bar{b}_1 is the m-multiple root. Therefore, the multiple poles are completely included in the solutions of equation (10). In the same way, we can assume $\mu_n = 1$.

2.4 The infinite poles case

In this subsection we will discuss the case in which polynomials are included. Set

$$f(z) = \sum_{k=0}^{m-1} c_k z^k + \sum_{l=m+1}^n \frac{\lambda_l}{1 - \bar{b}_l z},$$
 (21)

where c_{m-1} and λ_l (l = m + 1, ..., n) are nonzero. For the monomial z^{m-1} , when applying the backward shift operator to it, we have

$$\mathbf{S}(z^{m-1}) = \frac{z^{m-1} - 0}{z} = z^{m-2}.$$
 (22)

After mth backward shift operator taken on it, it becomes zero, viz,

$$\mathbf{S}^m(z^{m-1}) = 0. (23)$$

Then, for f(z) given by (21), we have

$$\begin{cases} \mathbf{S}(f)(z) = \sum_{k=1}^{m-1} c_k z^{k-1} + \sum_{l=1}^n \frac{\lambda_k \overline{b}_k}{1 - \overline{b}_k z} \\ \vdots \\ \mathbf{S}^{m-1}(f)(z) = c_{m-1} + \sum_{l=1}^n \frac{\lambda_k \overline{b}_k^{m-1}}{1 - \overline{b}_k z}, \\ \mathbf{S}^m(f)(z) = 0 + \sum_{l=1}^n \frac{\lambda_k \overline{b}_k^m}{1 - \overline{b}_k z}, \\ \vdots \\ \mathbf{S}^{m+n}(f)(z) = 0 + \sum_{l=1}^n \frac{\lambda_k \overline{b}_k^{m+n}}{1 - \overline{b}_k z}, \end{cases}$$

Similarly, because of the linear independence of $\{z^k\}$ and $\{\frac{1}{1-\overline{b}_l z}\}$, there exists a unique sequence $\{\mu_k\}_{k=0}^n$ such that

$$\mu_0 f(z) + \mu_1 \mathbf{S}(f)(z) + \dots + \mu_n \mathbf{S}^n(f)(z) = 0, \qquad (24)$$

which gives

$$\begin{pmatrix}
\mu_0 c_0 + \mu_1 c_1 + \mu_2 c_2 + \dots + \mu_{m-1} c_{m-1} = 0, \\
\mu_0 c_1 + \mu_1 c_2 + \mu_2 c_3 + \dots + \mu_{m-2} c_{m-1} = 0, \\
\vdots \\
\mu_0 c_{m-3} + \mu_1 c_{m-2} + \mu_2 c_{m-1} = 0, \\
\mu_0 c_{m-2} + \mu_1 c_{m-1} = 0, \\
\mu_0 c_{m-1} = 0,
\end{pmatrix}$$
(25)

and

$$\begin{cases} \mu_0 + \mu_1 \overline{b}_1 + \mu_2 \overline{b}_1^2 + \dots \mu_n \overline{b}_1^n = 0, \\ \vdots \\ \mu_0 + \mu_1 \overline{b}_n + \mu_2 \overline{b}_n^2 + \dots \mu_n \overline{b}_n^n = 0. \end{cases}$$
(26)

Since $c_{m-1} \neq 0$, from equations (25) we can see $\mu_0 = \mu_1 = \dots = \mu_{m-1} = 0$. While from equations (26), by the same way, we can see that $\{\bar{b}_k\}_{k=m+1}^n$ are solutions of equation

$$\mu_0 + \mu_1 x + \dots + \mu_{n-1} x^{n-1} + \mu_n x^n = 0, \qquad (27)$$

which can be simplified to

$$x^{m}(\mu_{n}x^{n-m} + \mu_{n-1}x^{n-m-1}... + \mu_{m}) = 0.$$
 (28)

It is easy to find that the truth $\{\bar{b}_k\}_{k=m+1}^n$ are exactly the solutions of

$$\mu_n x^{n-m} + \mu_{n-1} x^{n-1} \dots + \mu_m = 0$$

Without loss of generality, we can assume $\mu_n = 1$.

For the multiple $\{\overline{b}_k\}$ case, the discussion is similar to the above.

2.5 Sum-up

For any n-tuple $\{a_1, ..., a_n\}$ in \mathbb{D} , the corresponding finite orthogonal system

$$B_{\{a_1\}}, B_{\{a_1,a_2\}}, ..., B_{\{a_1,a_2,...,a_n\}},$$

where $B_{\{a_1,...,a_k\}}$ (k = 1,...,n) is defined by (1), is called an *n*-Blaschke system. It is proven in [22] that the *n*-Blaschke system may be obtained through Gram-Schmidt orthogonalization process from the reproductive kernels of order n,

$$e_{\{a_1\}}, e_{\{a_2\}}, ..., e_{\{a_n\}},$$

where if $a_k \neq 0$ has multiplicity l in the *n*-tuple $(a_1, ..., a_n)$, then

$$e_{\{a_k\}}(z) = \frac{1}{(1 - \overline{a}_k z)^j} (j = 1, ..., l);$$

if $a_k = 0$ has multiplicity l, then

$$e_{\{a_k\}}(z) = z^{j-1}(j = 1, ..., l).$$

Any rational function consists of the kernels can be represented by linear combination of the corresponding finite orthogonal system. The studies in the above three subsections through backward shift operator perfectly correspond to the cases of the poles of the rational functions. Our conclusion is that:

Theorem 2 Poles (including the zero and multiplicities) of rational functions are exactly solutions of an algebraic equation which can be derived through applying the backward shift operator.

3 Backward shift algorithm

With Theorem 2, we can estimate poles of LTI systems based on frequency responses. Generally, the measured data are corrupted by noise and errors of measurement. These will lead to some errors in estimating coefficients $\{\mu_k\}_{k=0}^{n-1}$ and finally affect the accuracy of the estimated poles.

It is here assumed a set of measurements $\{E_k\}_{k=1}^N$ from a system f(z), f(z) be a proper rational function, is given by

$$E_k = f(e^{i\omega_k}) + v_k, \tag{29}$$

where $\omega_k = \frac{2\pi(k-1)}{N}$, k = 1, ..., N, and v_k is the corruption to $f(e^{i\omega_k})$. The corruption $\{v_k\}$ can be used to model a number of different error sources. It can be assumed to either be a bounded deterministic sequence, or a stochastic process with zero-mean and bounded covariance.

Without losing the accuracy, we consider the equation (9) in the form

$$\mu_0 f(z) + \mu_1 \mathbf{S}(f)(z) + \dots + \mu_{n-1} \mathbf{S}^{n-1}(f)(z) = \mathbf{S}^n(f)(z).$$
(30)

The corresponding algebraic equation to (10) becomes

$$x^{n} - \mu_{n-1}x^{n-1} - \mu_{n-2}x^{n-2} - \dots - \mu_{1}x - \mu_{0} = 0.$$
(31)

We call our algorithm *backward shift algorithm* given as follows.

Algorithm 1 It has four steps:

- Generates data sets by taking backward shift operator to {E_k}^N_{k=1} n times.
- Determine the coefficients μ_k 's by relation (30).
- Find solutions of the algebraic equation (31) which gives the poles of the rational function.
- Find the composition coefficients by using a least-squares criterion.

The following are processes in detail. Firstly, based on $\{f(e^{i\omega_k})\}_{k=1}^N$, $\{S(f)(e^{i\omega_k})\}_{k=1}^N$ can be obtained by

$$\mathbf{S}(f)(e^{i\omega_k}) = \frac{f(e^{i\omega_k}) - f(0)}{e^{i\omega_k}}.$$
 (32)

With f(0) approximated by

$$f(0) = \frac{1}{2\pi} \int_0^{2\pi} f(e^{i\omega}) d\omega \approx \frac{1}{N} \sum_{k=1}^N f(e^{i\omega_k}),$$

we denote

$$\widetilde{\mathbf{S}}(f)(e^{i\omega}) = \frac{f(e^{i\omega_k}) - \frac{1}{N} \sum_{k=1}^N f(e^{i\omega_k})}{e^{i\omega_k}}.$$

Repeating the process, we have

$$\mathbf{S}^{l-1}f(0) = \frac{1}{2\pi} \int_0^{2\pi} \mathbf{S}^{l-1}f(e^{i\omega})d\omega$$
$$\approx \frac{1}{N} \sum_{k=1}^N \widetilde{\mathbf{S}}^{l-1}f(e^{i\omega_k}),$$

and

$$\widetilde{\mathbf{S}}^{l}(f)(e^{i\omega}) = \frac{\widetilde{\mathbf{S}}^{l-1}(f)(e^{i\omega}) - \frac{1}{N} \sum_{k=1}^{N} \widetilde{\mathbf{S}}^{l-1}f(e^{i\omega_{k}})}{e^{i\omega}}$$
$$\approx \mathbf{S}^{l}(f)(e^{i\omega}), \qquad (33)$$

The approximating data sets $\{\widetilde{\mathbf{S}}^2(f)(e^{i\omega_k})\}, \{\widetilde{\mathbf{S}}^3(f)(e^{i\omega_k})\}, \dots, \{\widetilde{\mathbf{S}}^n(f)(e^{i\omega_k})\}\$ are obtained. Therefore we obtain a linear system, denoted by

$$A\mu_N = b, \tag{34}$$

where $A = (A_{pq})_{N \times n}$, $A_{pq} = \widetilde{S}^{q-1}(f)(e^{i\omega_p})$, and μ , b are the columns given by

$$\mu = [\mu_0 \ \mu_1 \ \dots \ \mu_{n-1}]^T, b = [\widetilde{S}^n(f)(e^{i\omega_1}) \ \widetilde{S}^n(f)(e^{i\omega_2}) \ \dots \ \widetilde{S}^n(f)(e^{i\omega_N})]^T,$$

T being the transposition.

The problem (34) is a typical least-squares problem, since f(z), $\mathbf{S}(f)(z)$,..., $\mathbf{S}^{n-1}(f)(z)$ are independent. To solve (34) is equivalent to solve

 $\mu_N =$

$$\arg\min_{\mu} \frac{1}{N} \sum_{k=1}^{N} |\sum_{l=1}^{n} \mu_{l-1} \widetilde{\mathbf{S}}^{l-1}(f)(e^{i\omega_k}) - \widetilde{\mathbf{S}}^n(f)(e^{i\omega_k})|^2,$$
(35)

whose normal equations are

$$R_N \mu_N = \frac{1}{N} A^T b,$$

where $R_N = \frac{1}{N} A^T A$,

$$(R_N)_{pq} = \frac{1}{N} \sum_{k=1}^{N} \widetilde{\mathbf{S}}^{p-1}(f)(e^{i\omega_k}) \widetilde{\mathbf{S}}^{q-1}(f)(e^{i\omega_k}), \quad (36)$$

p, q = 1, 2, ..., n.

We have the following result.

Lemma 3 For the matrix R_N defined by (36), there exists a matrix R^* , such that

$$R^* = \lim_{N \to \infty} R_N. \tag{37}$$

Proof. It is easy to prove the result since $\mathbf{S}^{p-1}(f)(z)$ (p = 1, ..., n) are rational. First we can see

$$\lim_{N \to \infty} \widetilde{\mathbf{S}}(f)(e^{i\omega_k}) = \lim_{N \to \infty} \frac{f(e^{i\omega_k}) - \frac{1}{N} \sum_{l=1}^N f(e^{i\omega_l})}{e^{i\omega_k}}$$
$$= \frac{f(e^{i\omega_k}) - \frac{1}{2\pi} \int_0^{2\pi} f(e^{i\omega}) d\omega}{e^{i\omega_k}}$$
$$= \frac{f(e^{i\omega_k}) - f(0)}{e^{i\omega_k}}$$
$$= \mathbf{S}(f)(e^{i\omega_k}).$$
(38)

Therefore, according to (38), we have

$$\lim_{N \to \infty} \widetilde{\mathbf{S}}^2(f)(e^{i\omega_k}) = \frac{\mathbf{S}(f)(e^{i\omega_k}) - \mathbf{S}(f)(0)}{e^{i\omega_k}}$$
$$= \mathbf{S}^2(f)(e^{i\omega_k}).$$

Likewise, uniformly in k,

$$\lim_{N \to \infty} \widetilde{\mathbf{S}}^{p-1}(f)(e^{i\omega_k}) = \mathbf{S}^{p-1}(f)(e^{i\omega_k}).$$

It means for each element of R_N ,

$$\lim_{N \to \infty} (R_N)_{pq}$$

= $\lim_{N \to \infty} \frac{1}{N} \sum_{k=1}^{N} \widetilde{\mathbf{S}}^{p-1}(f)(e^{i\omega_k}) \widetilde{\mathbf{S}}^{q-1}(f)(e^{i\omega_k})$
= $\frac{1}{2\pi} \int_0^{2\pi} \mathbf{S}^{p-1}(f)(e^{i\omega}) \mathbf{S}^{q-1}(f)(e^{i\omega}) d\omega.$

Accordingly, R_N converges when $N \to \infty$. Denote

$$\lim_{N \to \infty} R_N = R^*.$$

The proof is complete.

The rank of R^* is n, because $f, \mathbf{S}(f), ..., \mathbf{S}^{n-1}(f)$ are linear independent. Similarly, it can be proved there exists a b^* given by

$$b^* = \lim_{N \to \infty} \frac{1}{N} Ab.$$

From the proof of lemma 3, we notice if μ^* is the true solution of (30), then there holds

$$\lim_{N \to \infty} \mu_N = \mu^*. \tag{39}$$

In the noise case, assume the noise $\{v_k\}$ is bounded. Denote $\{\widetilde{\mathbf{S}}_d(f)(e^{i\omega_k})\}, \{\widetilde{\mathbf{S}}_d^2(f)(e^{i\omega_k})\}, \{\widetilde{\mathbf{S}}_d^3(f)(e^{i\omega_k})\}, \dots, \{\widetilde{\mathbf{S}}_d^n(f)(e^{i\omega_k})\}\)$ as the data sets obtained through the backward shift processes based on $\{E_k\}_{k=1}^N$, *i.e.* corrupted by noise. The following result states a property of this algorithm.

Theorem 4 Suppose the noise $\{v_k\}$ is bounded by $\epsilon > 0$, that is, $|v_k| < \epsilon$. Let μ^* be the true solution of (30), $R^*\mu^* = b^*$, and μ'_N be the solution given by (35) with the sequence of data $\{\{\widetilde{\mathbf{S}}_d^l(f)(e^{i\omega_k})\}_{k=1}^N\}_{l=1}^n$ then there holds

$$\lim_{\substack{\epsilon \to 0, \\ N \to \infty}} \mu'_N = \mu^*.$$
(40)

Proof.

By substituting the measured data $\{E_k\}_{k=1}^N$ in the backward shift processes, the first data set is

$$\begin{split} \widetilde{\mathbf{S}}_{d}(f)(e^{i\omega_{k}}) \\ &= \frac{E_{k} - \frac{1}{N} \sum_{l=1}^{N} E_{l}}{e^{i\omega_{k}}} \\ &= \frac{f(e^{i\omega_{k}}) - \frac{1}{N} \sum_{l=1}^{N} f(e^{i\omega_{l}})}{e^{i\omega_{k}}} + \frac{v_{k} - \frac{1}{N} \sum_{l=1}^{N} v_{l}}{e^{i\omega_{k}}} \\ &= \widetilde{\mathbf{S}}(f)(e^{i\omega_{k}}) + \widetilde{\mathbf{S}}E(v_{k}), \end{split}$$

where $\widetilde{\mathbf{S}}E(v_k)$ stands for the error between $\widetilde{\mathbf{S}}_d(f)(e^{i\omega_k})$ and $\widetilde{\mathbf{S}}(f)(e^{i\omega_k})$. We can see

$$\lim_{\substack{\epsilon \to 0, \\ N \to \infty}} \widetilde{\mathbf{S}} E(v_k) = \lim_{\substack{\epsilon \to 0, \\ N \to \infty}} \frac{v_k - \frac{1}{N} \sum_{l=1}^N v_l}{e^{i\omega_k}} = 0.$$
(41)

The other data sets can be obtained from the following recursive formula, for m = 2, 3, ..., n,

$$\widetilde{\mathbf{S}}_{d}^{m}(f)(e^{i\omega_{k}}) = \widetilde{\mathbf{S}}^{m}(f)(e^{i\omega_{k}}) + \widetilde{\mathbf{S}}^{m}E(v_{k}), \qquad (42)$$

where

$$\widetilde{\mathbf{S}}^{m}(f)(e^{i\omega_{k}}) = \frac{\widetilde{\mathbf{S}}^{m-1}(f)(e^{i\omega_{k}}) - \frac{1}{N}\sum_{l=1}^{N}\widetilde{\mathbf{S}}^{m-1}(f)(e^{i\omega_{l}})}{e^{j\omega_{k}}},$$
$$\widetilde{\mathbf{S}}^{m}E(v_{k}) = \frac{\widetilde{\mathbf{S}}^{m-1}E(v_{k}) - \frac{1}{N}\sum_{l=1}^{N}\widetilde{\mathbf{S}}^{m-1}E(v_{l})}{e^{i\omega_{k}}}.$$
(43)

The error between $\widetilde{\mathbf{S}}_{d}^{m}(f)(e^{i\omega_{k}})$ and $\widetilde{\mathbf{S}}^{m}(f)(e^{i\omega_{k}})$ is $\widetilde{\mathbf{S}}^{m}E(v_{k})$. From (41), by mathematical induction, we also have

$$\lim_{\substack{\epsilon \to 0, \\ N \to \infty}} \widetilde{\mathbf{S}}^m E(v_k) = 0.$$
(44)

The equations for μ are

$$(A + \delta_A)\mu = b + \delta_b, \tag{45}$$

where $A_{pq} = \widetilde{\mathbf{S}}^{q-1}(f)(e^{i\omega_p}), \ \delta_{A_{pq}} = \widetilde{\mathbf{S}}^{q-1}E(v_p)(p,q = 1,2,...,n)$ and

$$b = [\widetilde{\mathbf{S}}^{n}(f)(e^{i\omega_{1}}) \ \widetilde{\mathbf{S}}^{n}(f)(e^{i\omega_{2}}) \ \dots \ \widetilde{\mathbf{S}}^{n}(f)(e^{i\omega_{N}})]^{T},$$

$$\delta_{b} = [\widetilde{\mathbf{S}}^{n}E(v_{1}) \ \widetilde{\mathbf{S}}^{n}E(v_{2}) \ \dots \ \widetilde{\mathbf{S}}^{n}E(v_{N})]^{T}.$$

The normal equations become

$$(R_N + \delta_{R_N})\mu = \overline{b} + \delta_{\overline{b}},\tag{46}$$

where

$$R_N = \frac{1}{N} A^T A,$$

$$\delta_{R_N} = \frac{1}{N} (A^t \delta_A + (\delta_A)^T A + (\delta_A)^T \delta_A),$$

$$\bar{b} = \frac{1}{N} A^T b,$$

$$\delta_{\bar{b}} = \frac{1}{N} ((\delta_A)^T b + (\delta_A)^T \delta_b + A \delta_b).$$

Since each element of δ_A and δ_b tends to 0 as $N \to \infty$ and $\epsilon \to 0$, besides each element of matrices A and b converges, then from (41) and (44), we can see that

$$\lim_{\substack{\epsilon \to 0, \\ N \to \infty}} \delta_{R_N} = 0, \quad \lim_{\substack{\epsilon \to 0, \\ N \to \infty}} \delta_{\overline{b}} = 0.$$
(47)

Let μ_N be the solution using data without noise and μ'_N be solution with noise. That is,

$$R_N \mu_N = \overline{b},$$
$$(R_N + \delta_{R_N})\mu'_N = \overline{b} + \delta_{\overline{b}}.$$

By perturbation theory [24], there is a result

$$\|\mu_{N}'-\mu_{N}\| \leq \|\mu_{N}\| \frac{\kappa(R_{N})}{1-\kappa(R_{N})\frac{\|\delta_{R_{N}}\|}{\|R_{N}\|}} \left(\frac{\|\delta_{R_{N}}\|}{\|R_{N}\|} + \frac{\|\delta_{\overline{b}}\|}{\|\overline{b}\|}\right),$$
(48)

where $\kappa(R_N) = ||R_N|| ||R_N^{-1}||$ is the condition number of R_N . We have the following error from μ'_N to μ^*

$$\begin{aligned} \|\mu'_{N} - \mu^{*}\| &\leq \|\mu'_{N} - \mu_{N}\| + \|\mu_{N} - \mu^{*}\| \\ &\leq \|\mu_{N}\| \frac{\kappa(R_{N})}{1 - \kappa(R_{N}) \frac{\|\delta_{R_{N}}\|}{\|R_{N}\|}} \left(\frac{\|\delta_{R_{N}}\|}{\|R_{N}\|} + \frac{\|\delta_{\overline{b}}\|}{\|\overline{b}\|} \right) \\ &+ \|\mu_{N} - \mu^{*}\|, \end{aligned}$$
(49)

From (48), it can be seen that

$$\lim_{\substack{\epsilon \to 0, \\ N \to \infty}} \|\mu'_N - \mu_N\| = 0; \tag{50}$$

On the other hand, from (39), there holds

$$\lim_{N \to \infty} \|\mu_N - \mu^*\| = 0.$$
 (51)

Therefore, we have

$$\lim_{\substack{\epsilon \to 0, \\ N \to \infty}} \|\mu'_N - \mu^*\| = 0.$$
(52)

The proof is complete.

The inequality (49) also gives the error estimation in the matrix norm sense. Once μ'_N is obtained, the poles $\{\overline{a}_k\}$

can be acquired by solving equation (31) again. Using the least-squares criterion with $\{E_k\}_{k=1}^N$, the corresponding coefficients can be obtained.

Remark 5 For order estimation of a rational function, by perturbation theory of matrix, there is a minimal perturbation that affects the rank of a matrix. Assume that the transfer function is of type (8) or (11), it is possible to determine the order through backward shift operator. We can repeat the shift processes, study the changes of A's rank in (34) and stop until A's rank does not change any more.

Remark 6 Observation and discussion show that if the transfer function G(z) is not a rational function, then the proposed backward shift method for a prescribed n may have no relevance with the solution of the best approximation to G(z) by rational functions of order not larger than n ([2]), although it does give the unique solution of the n-best approximation problem, that is G(z) itself, when G(z) is a rational function of order n.

4 Numerical examples

In this section, we give some examples for illustrating the proposed backward shift algorithm, the first one is

$$\begin{split} f(z) &= \\ \frac{1}{1-0.5z} - \frac{3}{1+0.6z} + \frac{2}{1+0.7z} + \frac{3-5i}{6(1-(0.5+0.3i)z)} \\ &+ \frac{3+5i}{6(1-(0.5-0.3i)z)}, \end{split}$$

in which each pole is different.

Table 1

 $\{\overline{a}_k\}$ obtained in no noise case for example 1.

data number	\overline{a}_1	\overline{a}_2	\overline{a}_3
15	-0.7068	-0.5854	0.4712 - 0.2740i
20	-0.7005	-0.5992	0.4975 - 0.2999i
25	-0.7	-0.6	0.4999 + 0.3i
28	-0.7	-0.6	0.5 + 0.3i

Table 2

Continued $\{\overline{a}_k\}$ in no noise case for example 1.

data number	\overline{a}_4	\overline{a}_5
15	0.4712 + 0.2740i	0.2623
20	0.4975 + 0.2999i	0.4893
25	0.4999 - 0.3i	0.4997
28	0.5 - 0.3i	0.5

Table 3

 $\{\overline{a}_k\}$ obtained for example 1 with noised data(m = 20).

	SNR	\overline{a}_1	\overline{a}_2	\overline{a}_3
	10	-0.8194 - 0.0971i	0.5130 - 0.3587i	-0.3527 + 0.1973i
Ì	20	-0.6620 - 0.0454i	-0.6411 + 0.0888i	0.4526 + 0.2298i
	30	-0.7195 + 0.0086i	-0.5355 - 0.0331i	0.4767 - 0.2672i
	40	-0.7208 + 0.0006i	-0.5775 - 0.0001i	0.4874 - 0.2837i

Table 4 Continued $\{\overline{a}_k\}$ obtained for example 1 with noised data(m = 20).

SNR	\overline{a}_4	\overline{a}_5
10	0.4443 + 0.1699i	0.1049 + 0.3324i
20	0.4767 - 0.3016i	0.2041 - 0.2514i
30	0.4532 + 0.2650i	0.0802 + 0.1173i
40	0.4555 + 0.2930i	0.3667 + 0.0812i



Fig. 1. Comparison of different noised data with SNR=10,20,30,40, respectively, for example 1. And the data number m=20.

Table 1 and 2 show the results in no noise case, we can locate true poles of original system by using no more than 30 frequency responses for example 1. Table 3 and 4 are results by using the noised data (m=20) in each SNR(signal-to-noise ratio) level. Figure 1 shows frequency responses of approximating systems obtained by using noised data in different levels of SNR. We can see the approximations with order 5 is very good.

The second example is

$$f(z) = \frac{1}{1 - 0.9048z} + \frac{z^2 - 1}{(1 - 0.3679z)^2},$$
 (53)

in which 0.3679 has multiplicity 2.

Table 5

 $\{\overline{a}_k\}$ obtained in no noise case for example 2.

data number	\overline{a}_1	\overline{a}_2	\overline{a}_3
m=50	0.9040	0.3655 + 0.0346i	0.3655 - 0.0346i
m=80	0.9048	0.3679 - 0.0017i	0.3679 + 0.0017i
m=120	0.9048	0.3679	0.3679

Table 5 shows in the non-noise case, true solution can be gained. Table 6 are the results using the noised data. We can see the estimated poles are close to the original ones. In Figure 2, it shows the frequency responses of approximating systems obtained by using data added up with Gaussian noise in different levels of SNR.

Table 6						
$\{\overline{a}_k\}$ obtained	\mathbf{for}	example	2	with	noised	data.

\mathbf{SNR}	\overline{a}_1	\overline{a}_2	\overline{a}_3
10	0.8627 + 0.0008i	-0.4135 + 0.0591i	-0.0218 - 0.0374i
20	0.8657	-0.0898 - 0.1313i	-0.0842 + 0.1163i
30	0.8800	0.2316 + 0.2253i	0.2229 - 0.2176i
40	0.8999	0.3537 + 0.0866i	0.3502 - 0.0865i



Fig. 2. Approximations of different noised data with SNR= 10,20,30,40, respectively, for example 2.

In the third example, we compare our algorithm with the existing algorithms given by [8] and [13]. We choose a simple second order system for comparison. Consider

$$g(z) = \frac{4}{z - 0.5} + \frac{5}{z - 0.2},\tag{54}$$

for algorithms in [8] and [13], it is to estimate poles 0.5 and 0.2 with measured data. The function g(z) is analytic outside of the unit disc, by taking z = 1/z, then it is analytic in the unit disc and the proper part is

$$f(z) = \frac{4}{1 - 0.5z} + \frac{5}{1 - 0.2z},\tag{55}$$

we are to estimate the points 0.5 and 0.2 of f(z) by the proposed method. The results are shown in the table 7 and 8 for both noise and no noise cases, respectively. All the data are equally spaced in the corresponding interval. In the noise case, the data are added up with Gaussian noise at SNR = 20. We can see our results are better.

Table 7						
Results of a_1	and a_2	with :	no noise	for	example 3.	

rr						
methods	data a1 10 0.4999		a_2			
Proposed			0.1997			
Proposed	15	0.5	0.2			
Nara-Ando's [13]	a-Ando's [13] 10 -0.9		0.9646 - 0.0314i			
Kang-Lee's [8]	30	-0.9638 - 1.0765i	-0.3846 + 0.8661i			

Table 8 Results of a_1 and a_2 with noise for example 3.

methods	data	a_1	a_2
Proposed	20	0.4935 + 0.0196i	0.2080 + 0.0193i
Nara-Ando's [13]	50	0.8888 - 0.4782i	-0.6116 + 0.6413i
Kang-Lee's [8]	50	0.5384 + 0.6089i	0.0557 - 0.0774i

5 Conclusions

This paper presents a new method for estimating poles of LTI systems by using backward shift operator. The numerical results show this algorithm is practical. The number of used data is not large. In fact, for the conjugate symmetry of the values on the unit circle of the transfer function, it only needs data that in the interval $(0, \pi)$, the rest on the other half circle can be denoted to be the conjugated value of the measured data, which decreases the cost of experiments. In each shift process, it computes a mean of errors, this will contribute to eliminate the large errors. From the illustrating examples, we can see this algorithm is convenient and efficient.

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