Recursive Identification for Hammerstein Systems with State-space Model

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In this paper, the subspace methods for multi-input multi-output (MIMO) Hammerstein systems in state-space form Abstract are considered, for the reason that relation of input and output in MIMO case is appropriate to indicate in this form. In open-loop case, a subspace identification method is given for Hammerstein systems, in which the nonlinear static function can be represented as a linear combination of finite basis functions, and the recursive version of this method is also given. We show that, in mild conditions, the method given in this paper is consistent in some sense. A numerical example is given to show the performance of this method. Key words Subspace method, Hammerstein system, recursive method, principal component analysis (PCA), multi-input multioutput (MIMO), stochastic approximation

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Generally, system identification is carried out to figure out the structure of the dynamic system by making use of the system input and output corrupted by additive noise. Most of the existing methods focus on the parametric model-based techniques to find a well-suited model described by a finite number of data. Examples of such methods include the prediction error method (PEM) to estimate the free parameters based on an optimization problem, the instrumental variable (IV) method to find a good principal random variable uncorrelated with the noise and sufficiently correlated with the states and inputs, and the maximum likelihood method and the subspace identification method $(SIM)^{[1-4]}$, and so on, respectively, in various model structure.

The subspace identification method is designed especially for the time-invariant systems, the beginning of which could be traced back to the 1960s. It is naturally applicable to the multi-input multi-output (MIMO) systems in the form of state space model^[5], which is convenient for estimation, filtering, prediction, and control. Much attention has been attracted from researchers because of its efficient, convenient, and stable computational characteristics. As the subspace method has been developed in several decades, it has attained maturity, especially against linear state space model.

As for the nonlinear system, the state space model description is applied widely in various fields. General form of the dynamic system could be described as follows:

$$x_{t+1} = f(t, x_t, u_t, \theta_t) + v_{t+1}$$
(1)

$$y_t = h(t, x_t, u_t, \theta_t) + w_t \tag{2}$$

in which x_t , u_t , and y_t are the system state, input, and output respectively, and θ_t denotes the parameter of the system, v_t and w_t represent the process noise and the observed noise.

In the nonlinear system field, there is a simple and typical kind of system, which consists of a linear dynamic part cascaded with a static nonlinear part at the input or output. The system is called Hammerstein system, if the nonlinearity is prior to the linear part. And the same is called Wiener system, if the cascade order alters. The Hammerstein system is used for modeling biological process, chemical process, and in signal processing applications, which can be shown in the state-space form as

$$x_{k+1} = Ax_k + Bg(u_k) + v_{k+1} \tag{3}$$

$$y_{k+1} = Cx_{k+1} + Du_{k+1} + w_{k+1} \tag{4}$$

Nonlinear identification^[6-7] is an increasingly active research topic, which includes parametric and nonparametric methods. As is widely known, no omnipotent identification method exists. Extensions to nonlinear systems though desirable are much harder in general. There are plenty of methods for Hammerstein system identification^[8-15], but they are mainly designed against the single-input and single-output (SISO) $model^{[16-17]}$. These days, researches on the subspace methods have been extended from linear systems to nonlinear systems. And some results of the nonlinear system identification methods based on subspace identification have appeared in succession^[18-21]</sup>. However, the nonlinear functions of these existing subspace identification methods are in particular forms and little theoretical analysis about convergence, consistency [8-9, 22-23], rate, and so on, has been given.

This paper is targeted to give a subspace identification method of the Hammerstein system in open-loop case, based on the state space representation. And the recursive algorithm $^{[24-26]}$ of the method is also shown. In order to obtain faster computation, we have utilized the averaging technique in view of the principal analysis based on stochastic approximation [27-28] to identify the extended observable matrix. In addition, some asymptotic quality is analyzed under mild condition. And a simulated example is given to show that the average technique is indeed a favorable alternative to speed up the computation process.

The rest of this paper is organized as follows: Section 1 sets the model and its assumptions, which hold throughout the whole paper. Section 2 gives the recursive version of identification, and its consistency analysis is also followed. Finally, a simulated example is given in Section 3 to evaluate the given method.

Notations.

E : Expectation;

R : Real field:

 \mathbf{R}^n : Real space in *n* dimensions;

 A^{T} : Transpose of A;

 $\lambda_{\max}(A)$: Largest eigenvalue of A;

 $\operatorname{span}(A)$: Space spanned by column of A;

 $\operatorname{Adj}(A)$: Adjoint matrix of A.

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1.1 Model and assumptions

We consider the discrete MIMO Hammerstein system in state space form as follows:

$$\boldsymbol{x}_{k+1} = A\boldsymbol{x}_k + B\boldsymbol{h}(\boldsymbol{u}_k) \tag{5}$$

$$\boldsymbol{y}_{k+1} = C \boldsymbol{x}_{k+1} + \boldsymbol{e}_{k+1}$$
 (6)

where $\boldsymbol{x}_k \in \mathbf{R}^n$, $\boldsymbol{u}_k \in \mathbf{R}^m$, $\boldsymbol{y}_k \in \mathbf{R}^l$, $\boldsymbol{e}_k \in \mathbf{R}^l$ are the system state, input, measured output, and observed noise, respectively. Here, $A \in \mathbf{R}^{n \times n}$, $B \in \mathbf{R}^{n \times m}$, $C \in \mathbf{R}^{l \times n}$, and $\boldsymbol{h} : \mathbf{R}^m \to \mathbf{R}^m$ is the unknown static nonlinear function. We assume the order of the system has been known beforehand, denoted as n.

For system (5) and (6), we make the following assumptions:

A 1. (A, C) is observable and (A, B) is controllable;

A 2. The system (5) is asymptotically stable, i.e., $|\lambda_{\max}(A)| < 1;$

A 3. The input u_k is a sequence of independently and identically distributed random variables with m dimensions;

A 4. Let $\boldsymbol{\varphi}_i(\boldsymbol{u}_k)$, $i = 1, \dots, N$ be N bounded functions of m dimensions, satisfying $\mathrm{E}\{\boldsymbol{\varphi}_i(\boldsymbol{u}_k)\boldsymbol{\varphi}_j(\boldsymbol{u}_k)^{\mathrm{T}}\} = \delta_{ij}$, where

$$\delta_{ij} = \begin{cases} I, & i = j\\ 0, & i \neq j \end{cases}$$

and \boldsymbol{h} could be constructed as the sum of $\boldsymbol{\varphi}_i(\boldsymbol{u}_k)$ with proper coefficients, namely,

$$\boldsymbol{h}(\boldsymbol{u}) = \sum_{i=1}^{N} D_i \boldsymbol{\varphi}_i(\boldsymbol{u})$$

where $D_1 = I$;

A 5. $\{e_k\}$ is a sequence of bounded white noise, independent from $\{u_k\}$ satisfying

$$\mathrm{E}\{\boldsymbol{e}_k\} = \boldsymbol{0}, \ \mathrm{E}\{\boldsymbol{e}_k\boldsymbol{e}_k^{\mathrm{T}}\} = \sigma_e I$$

where $\sigma_e > 0$ and I is the unit matrix.

Remark 1. Since N could be taken large enough, A 4 is reasonably acceptable, if **h** has a wide range. Any nonlinear function could be approximated as $\mathbf{h}(\mathbf{u}) = \sum_{i=1}^{N} D_i \varphi_i(\mathbf{u})$. This condition makes sure that $\frac{1}{L} \sum_{k=l}^{L} E\{\mathbf{h}(\mathbf{u}_k)\mathbf{h}(\mathbf{u}_{k-l})^{\mathrm{T}}\} = \sum_{i=1}^{N} D_i D_i^{\mathrm{T}} = I + \sum_{i=2}^{N} D_i D_i^{\mathrm{T}} > 0$, which can be considered as some kind of persistently excited restriction. Furthermore, the limitation on $D_1 = I$ could be broadened to nonsingular matrix. In this case, we could multiply it to B from the right, changing B into BD, and $\mathbf{h}(\mathbf{u})$ could still be approximated by a sum of the basic functions with proper coefficients and I as the first coefficient.

Remark 2. The problem of approximation error of the nonlinear function \boldsymbol{h} on the result of the estimation is $O(\epsilon)$, which can be obtained from the theory of robustness of stochastic approximation-based principal component analysis (SABPCA) algorithm later introduced. Certain conditions have been fulfilled in order to guarantee its estimation accuracy. And further knowledge on this topic can be found in [29].

1.2 Denotation

The purpose of identification is to estimate the system matrices A, B, C, and the coefficients $\{D_i, i = 2, \dots, N\}$ of the unknown nonlinear function \boldsymbol{h} .

In order to identify the extended state space model corresponding to the system (5) and (6), let $D = [I, D_2, \dots, D_N] \in \mathbf{R}^{m \times Nm}, \quad \tilde{\boldsymbol{u}}_k = [\boldsymbol{\varphi}_1^{\mathrm{T}}(\boldsymbol{u}_k), \dots, \boldsymbol{\varphi}_N^{\mathrm{T}}(\boldsymbol{u}_k)]^{\mathrm{T}} \in \mathbf{R}^{Nm}$, and then the state space model can be described as

$$\begin{aligned} \boldsymbol{x}_{k+1} &= A\boldsymbol{x}_k + BD\boldsymbol{u}_k \quad (7) \\ \boldsymbol{y}_{k+1} &= C\boldsymbol{x}_{k+1} + \boldsymbol{e}_{k+1} \quad (8) \end{aligned}$$

$$\boldsymbol{y}_{k+1} = C \boldsymbol{x}_{k+1} + \boldsymbol{e}_{k+1}$$
 (8)

Let

$$\boldsymbol{y}_{f}(k) = \begin{bmatrix} \boldsymbol{y}_{k-f+1}^{\mathrm{T}} & \boldsymbol{y}_{k-f+2}^{\mathrm{T}} & \cdots & \boldsymbol{y}_{k}^{\mathrm{T}} \end{bmatrix}^{\mathrm{T}} \in \mathbf{R}^{lf}$$

and we denote $\boldsymbol{e}_{f}(k)$ and $\boldsymbol{u}_{f}(k)$ in the same form as $\boldsymbol{y}_{f}(k)$, the components of which are respectively consisted of $\boldsymbol{e}_{k-f+1}, \cdots, \boldsymbol{e}_{k}$ and $\boldsymbol{u}_{k-f+1}, \cdots, \boldsymbol{u}_{k}$ in stead of $\boldsymbol{y}_{k-f+1}, \cdots, \boldsymbol{y}_{k}$.

Denote

$$Y_k = \begin{bmatrix} \boldsymbol{y}_f(1) & \boldsymbol{y}_f(2) & \cdots & \boldsymbol{y}_f(k) \end{bmatrix} \in \mathbf{R}^{lf \times k}$$

and Y_k is the Hankle matrix of all the output data before the time k. f is an arbitrary integer satisfying f > n.

1.3 Identification of O_f

Having denoted all the symbols we need, it is easy to gain the extended state space model by the state space system (7) and (8) as follows:

$$\boldsymbol{y}_f(k) = O_f \boldsymbol{x}_{k-f+1} + \Phi_f \boldsymbol{u}_f(k) + \boldsymbol{e}_f(k)$$
(9)

where

$$O_f = \begin{bmatrix} C \\ CA \\ \vdots \\ CA^{f-1} \end{bmatrix} \in \mathbf{R}^{lf \times n}$$

is an extended observable matrix of rank n, and

$$\Phi_f = \begin{bmatrix} 0 & \cdots & 0 & 0\\ CBD & \cdots & 0 & 0\\ \vdots & \ddots & \vdots & \vdots\\ CA^{f-2}BD & \cdots & CBD & 0 \end{bmatrix} \in \mathbf{R}^{lf \times Nmf} \quad (10)$$

is a block lower triangular Toeplitz matrix.

In order to estimate O_f and noticing (9), we intend to remove the input term $\boldsymbol{u}_f(k)$.

Lemma 1. If A 2 ~ A 5 hold, then $\{(\boldsymbol{y}_k^{\mathrm{T}}, \boldsymbol{u}_k^{\mathrm{T}})^{\mathrm{T}}\}$ is an asymptotically stable and ergodic sequence; and so is $(\boldsymbol{y}_f(k)^{\mathrm{T}}, \boldsymbol{u}_f(k)^{\mathrm{T}})^{\mathrm{T}}$.

From Lemma 1, by the ergodic characteristic of $\boldsymbol{u}_{f}(k)$ and $\boldsymbol{y}_{f}(k)$ we could conclude that:

$$\frac{1}{k} Y_k U_k^{\mathrm{T}} \xrightarrow[k \to \infty]{a.s.} R_{yu} = \mathrm{E}\{\boldsymbol{y}_f(k)\boldsymbol{u}_f^{\mathrm{T}}(k)\}$$
$$\frac{1}{k} U_k U_k^{\mathrm{T}} \xrightarrow[k \to \infty]{a.s.} R_{uu} = \mathrm{E}\{\boldsymbol{u}_f(k)\boldsymbol{u}_f^{\mathrm{T}}(k)\}$$

A 3 and A 4 assure $R_{uu} > 0$, so $U_k U_k^{\mathrm{T}}/k$ is reversible for large enough k, and $(U_k U_k^{\mathrm{T}}/k)^{-1} \xrightarrow[k \to \infty]{a.s.} R_{uu}^{-1}$.

Considering

$$\boldsymbol{z}_{f}(k+1) = \boldsymbol{y}_{f}(k+1) - R_{yu}R_{uu}^{-1}\boldsymbol{u}_{f}(k+1)$$
$$R = \mathbb{E}\{\boldsymbol{z}_{f}(k+1)\boldsymbol{z}_{f}^{\mathrm{T}}(k+1)\} = R_{yy} - R_{yu}R_{uu}^{-1}R_{uy} \quad (11)$$

and

$$egin{aligned} r_{xx} &= \mathrm{E}\{oldsymbol{x}_k oldsymbol{x}_k^{\mathrm{T}}\}\ r_{xu} &= \mathrm{E}\{oldsymbol{x}_{k-f+1} oldsymbol{u}_f^{\mathrm{T}}(k)\} \end{aligned}$$

by the definition (11) of R, the extended state space model (9) and the assumptions, we have

$$R = O_f (r_{xx} - r_{xu} R_{uu}^{-1} r_{ux}) O_f^{\rm T} + \sigma_e I$$
 (12)

Since the extended observable matrix O_f is full column rank (rank $O_f = n$), by (12), if

$$r_{xx} - r_{xu} R_{uu}^{-1} r_{ux} > 0 (13)$$

it is concluded that n of the eigenvalues of R are larger than σ_e , and the rest of lf - n ones exactly equal to σ_e . From the singular value decomposition, we know that under condition (13), R can be decomposed as follows:

$$R = \begin{bmatrix} R_s & R_n \end{bmatrix} \begin{bmatrix} \Lambda_s & 0 \\ 0 & \Lambda_n \end{bmatrix} \begin{bmatrix} R_s^{\mathrm{T}} \\ R_n^{\mathrm{T}} \end{bmatrix}$$
(14)

where $\Lambda_n = \sigma_e I$. R_s includes the *n* principal eigenvectors corresponding to the *n* principal eigenvalues and R_n contains lf - n eigenvectors corresponding to the eigenvalue σ_e . $[R_s \ R_n]$ is orthogonal matrix.

Comparing (12) and (14), the space spanned by the columns of O_f is identical to the space spanned by the columns of R_s , namely,

$$\operatorname{span}(O_f) = \operatorname{span}(R_s)$$

Therefore, the problem turns to be condition under which the inequality (13) will hold. And the following proposition will give an answer to this problem.

Proposition 1. If $A \ 1 \sim A \ 5$ hold, then

$$r_{xx} - r_{xu} R_{uu}^{-1} r_{ux} > 0$$

Proof. The proof is given in Appendix A. Proposition 1 assures $\operatorname{span}(O_f) = \operatorname{span}(R_s)$. Therefore, in order to estimate O_f , we just have to estimate the first n eigenvectors of R in (14).

$\mathbf{2}$ **Recursive algorithm**

In this section, we will give a version of recursive subspace identification method of Hammerstein system. First of all, a recursive estimation of O_f will be shown.

Estimation of O_f $\mathbf{2.1}$

In the previous section, we have shown that O_f can be estimated from the *n* principal eigenvectors $\boldsymbol{v}^{(i)}$ of *R*. And the principal eigenvectors can be calculated from the stochastic approximation-based principal component analysis algorithm. Here, to accelerate the computation speed, we will utilize the averaging technique of stochastic $approximation^{[16]}$.

Denote $\boldsymbol{v}^{(i)}$ as a unit eigenvector of R corresponding to the eigenvalue $\lambda^{(i)}$, $i = 1, \cdots, lf$. In recursively estimating $\boldsymbol{v}^{(i)}$ and $\lambda^{(i)}$, $i = 1, \dots, lf$, we will compute R, defined in (11). So, we perform the data projection as analyzed in Subsection 1.3:

$$\hat{\boldsymbol{z}}_{f}(k+1) = \\ \boldsymbol{y}_{f}(k+1) - R_{yu}(k+1)R_{uu}^{-1}(k+1)\boldsymbol{u}_{f}(k+1) \quad (15)$$

$$R_{yu}(k+1) = R_{yu}(k) - \frac{1}{k+1} \left[R_{yu}(k) - \boldsymbol{y}_f(k+1) \boldsymbol{u}_f^{\mathrm{T}}(k+1) \right]$$
$$R_{uu}^{-1}(k+1) = N(k+1) = N(k) + \frac{1}{k} \left[N(k) - (k+1) \frac{N(k)\boldsymbol{u}_f(k)\boldsymbol{u}_f(k)N(k)}{k+\boldsymbol{u}_f^{\mathrm{T}}(k)N(k)\boldsymbol{u}_f(k)} \right]$$

The initial value $R_{yu}(0) = 0$ and N(0) could be taken as arbitrary positive matrix.

After obtaining the data projection, we will calculate the component eigenvectors of $\hat{\boldsymbol{z}}_{f}(k+1)\hat{\boldsymbol{z}}_{f}^{\mathrm{T}}(k+1)$ by principal component analysis (PCA) based on the stochastic approximation.

$$\begin{split} \overset{1)}{\boldsymbol{v}} & \quad \tilde{\boldsymbol{v}}_{k+1}^{(1)} = \boldsymbol{v}_{k}^{(1)} + a_{k} \hat{\boldsymbol{z}}_{f}(k+1) \hat{\boldsymbol{z}}_{f}^{\mathrm{T}}(k+1) \boldsymbol{v}_{k}^{(1)} \\ & \quad \boldsymbol{v}_{k+1}^{(1)} = \frac{\tilde{\boldsymbol{v}}_{k+1}^{(1)}}{\|\tilde{\boldsymbol{v}}_{k+1}^{(1)}\|}, \quad \text{ if } \|\tilde{\boldsymbol{v}}_{k+1}^{(1)}\| \neq 0 \end{split}$$

If $\|\tilde{\boldsymbol{v}}_{k+1}^{(1)}\| = 0$, then reset $\boldsymbol{v}_k^{(1)}$ to be the unit vector. $\boldsymbol{v}_{k}^{(1)}$ should be taken as the one of the k-th R's unit characteristic vector estimation.

Denote $V_k^{(j)} = [\boldsymbol{v}_k^{(1)}, P_k^{(1)} \boldsymbol{v}_k^{(2)}, \cdots, P_k^{(j-1)} \boldsymbol{v}_k^{(j)}]$ be a matrix of $lf \times j$, where 2) Assume that $\boldsymbol{v}_k^{(i)}$, $i = 1, \cdots, j$ has been defined.

$$P_k^{(i)} = I - V_k^{(i)} V_k^{(i)\dagger}, \quad i = 1, 2, \cdots, j-1$$

 $V_{\scriptscriptstyle \! k}^{(j)\dagger}$ here represents the Moore-Penrose pseudoinverse of $V_k^{(j)}$. When k is large enough, $V_k^{(j)}$ is a full rank matrix of order $lf \times j (lf > j)$

$$V_k^{(j)\dagger} = (V_k^{(j)\mathrm{T}} V_k^{(j)})^{-1} V_k^{(j)\mathrm{T}}$$

Furthermore, we could estimate the j-th unit characteristic vector:

$$\begin{split} \tilde{\boldsymbol{v}}_{k+1}^{(j+1)} = & P_k^{(j)} \boldsymbol{v}_k^{(j+1)} + a_k P_k^{(j)} \hat{\boldsymbol{z}}_f(k+1) \hat{\boldsymbol{z}}_f^{\mathrm{T}}(k+1) P_k^{(j)} \boldsymbol{v}_k^{(j+1)} \\ \boldsymbol{v}_{k+1}^{(j+1)} = & \frac{\tilde{\boldsymbol{v}}_{k+1}^{(j+1)}}{\|\tilde{\boldsymbol{v}}_{k+1}^{(j+1)}\|}, \quad \text{if } \|\tilde{\boldsymbol{v}}_{k+1}^{(j+1)}\| \ge \epsilon \end{split}$$

where $0 < \epsilon < 1/4$, $j = 1, 2, \cdots, lf - 1$. If $\|\tilde{\boldsymbol{v}}_{k+1}^{(j+1)}\| < \epsilon$, then reset $\boldsymbol{v}_{k}^{(j+1)}$, so that $\|\boldsymbol{v}_{k}^{(j+1)}\| =$ 1 and $||P_k^{(j)} \boldsymbol{v}_k^{(j+1)}|| = 1.$

Then, we will give the algorithm as follows.

Step 0. Let $a_k = 1/k$ and $0 < M_k \longrightarrow \infty$ and take $\{\boldsymbol{v}_0^{(j)}, j = 1, \cdots, lf\}$ as arbitrary unit orthogonal group, $\lambda_0^{(j)} = 0, \ j = 1, \cdots, lf, \ \sigma_0 = 0, \ k = 0.$

Step 1. Data projection to calculate
$$\boldsymbol{z}_f(k)$$
 using (15):

$$\hat{\boldsymbol{x}}_{f}(k+1) = \\ \boldsymbol{y}_{f}(k+1) - R_{yu}(k+1)N(k+1)\boldsymbol{u}_{f}(k+1)$$

Step 2. Use SABPCA algorithm to calculate the principal characteristic vectors $\boldsymbol{v}_{k+1}^{(i)}, i = 1, \cdots, lf$ as shown above.

Step 3. Estimate the eigen-roots

$$\begin{split} \lambda_{k+1}^{(j)} &= [\lambda_k^{(j)} - a_k (\lambda_k^{(j)} - \| \hat{\boldsymbol{z}}_f^{\mathrm{T}}(k+1) \boldsymbol{v}_k^{(j)} \|^2)] \times \\ & \mathbf{1}_{[|\lambda_k^{(j)} - a_k (\lambda_k^{(j)} - \| \hat{\boldsymbol{z}}_f^{\mathrm{T}}(k+1) \boldsymbol{v}_k^{(j)} \|^2)| \le M_{\sigma_k}]}, \\ & j = 1, 2, \cdots, lf \end{split}$$

where

where

$$\sigma_{k+1} = \sigma_k + \mathbf{I}_{[|\lambda_k^{(j)} - a_k(\lambda_k^{(j)} - \|\hat{\boldsymbol{z}}_f^{\mathrm{T}}(k+1)\boldsymbol{v}_k^{(j)}\|^2)| > M_{\sigma_k}]}$$

Step 4. if $k < N_0$, let k = k + 1 and go to Step 1, else use fast average technique to proceed the step:

$$\begin{split} & \pmb{v}_{k+1}^{(j)} = \frac{1}{k+1} (\pmb{v}_{k}^{(j)} + k \times \pmb{v}_{k+1}^{(j)}) \\ & \lambda_{k+1}^{(j)} = \frac{1}{k+1} (\lambda_{k}^{(j)} + k \times \lambda_{k}^{(j)}) \end{split}$$

and perform normalization to $\boldsymbol{v}_{k}^{(j)}$, $j = 1, \cdots, lf$; **Step 5.** Rearrange $\{\lambda_{k+1}^{(j)}, \boldsymbol{v}_{k+1}^{(j)}\}$, so that $\lambda_{k+1}^{(1)} \geq 0$ $\lambda_{k+1}^{(2)} \ge \cdots \ge \lambda_{k+1}^{(lf)}$, and let

$$\hat{O}_f(k+1) = [\boldsymbol{v}_{k+1}^{(1)} \ \boldsymbol{v}_{k+1}^{(2)} \ \cdots \ \boldsymbol{v}_{k+1}^{(n)}]$$
(16)

Step 6. Take k = k + 1 and return to Step 1.

Remark 1. From the analysis and algorithm, $\lambda^{(j)}$ is constant when j > n, and $\lambda^{(n)} > \lambda^{(n+1)}$, so the order n of the system matrix can be estimated from $\lambda_k^{(j)}$ if the order is unknown.

Remark 2. Averaging technique will make the noise converge in equilibrium, however, the state deviation in the beginning is too large to use averaging technique, since it will eliminate the deviation too slowly. Therefore, in the algorithm above, we do not utilize fast average technique in the first N_0 steps. And N_0 will be taken large enough so that the averaging technique plays its role more efficiently.

The following theorem indicates O_f to be strongly consistent up to some extent.

Theorem 1. If A 1 \sim A 5 hold, then the estimation $\hat{O}_f(k)$ obtained as in (16), converges to O_f almost everywhere, i.e., there exists a nonsingular matrix T of order $n \times n$, so that

$$\hat{O}_f(k) \xrightarrow[k \to \infty]{\text{a.s.}} O_f T$$

Proof. The proof of Theorem 1 is also attached in Appendix B.

2.2 Recursive estimation of A, B, C, and D

2.2.1Estimation of A and C

Let $\hat{A}(k)$, $\hat{B}(k)$, $\hat{C}(k)$, and $\hat{D}(k)$ be the k-th estimation of A, B, C, and D, respectively. Once O_f has been known, the system matrix C can be directly observed from O_f .

As for the estimation of A, we will utilize the transfer invariance of O_f :

$$O_f(l+1:lf,:) = O_f(1:lf-l,:)A$$

Denote

$$\underline{O}(k) = O_f(k)((l+1):lf,:)$$

$$\overline{O}(k) = \hat{O}_f(k)(1:(f-1)l,:)$$

utilizing the transfer invariance of O_f , then

$$\underline{O}(k+1) = O(k) \times A + \epsilon_k$$

where ϵ_k is the noise in the process of estimation.

Using recursive least square estimation method, $\hat{A}(k)$ could be recursively calculated as follows:

$$\begin{cases} \hat{A}(k+1) = \hat{A}(k) + \gamma(k)G(k)\overline{O}(k)[\underline{O}(k+1) - \overline{O}(k) \times \hat{A}(k)] \\ G(k) = G(k) - \gamma(k)G(k)\overline{O}(k)\overline{O}(k)^{\mathrm{T}}G(K) \end{cases}$$
(17)

where

$$G(k) = \left(\sum_{i=0}^{k-1} \overline{O}(i)\overline{O}(i)^{\mathrm{T}}\right)^{-1}$$
$$\gamma(k) = (1 + \overline{O}(k)^{\mathrm{T}}G(k)\overline{O}(k))^{-1}$$

2.2.2 Estimation of B and D

Let $O_f^{\perp} \in \mathbf{R}^{lf \times (lf-n)}$ be denoted as the arbitrary orthogonal complement of O_f , i.e. $O_f^{\mathrm{T}} O_f^{\perp} = 0$ and O_f^{\perp} is full rank in column.

To estimate matrix B and D, we pre-multiply (9) by $(O_f^{\perp})^{\mathrm{T}}$, then

$$(O_f^{\perp})^{\mathrm{T}} \boldsymbol{y}_f(k) = (O_f^{\perp})^{\mathrm{T}} \Phi_f \boldsymbol{u}_f(k) + (O_f^{\perp})^{\mathrm{T}} \boldsymbol{e}_f(k)$$
(18)

in which $\boldsymbol{u}_{f}(k)$ and $\boldsymbol{e}_{f}(k)$ are uncorrelated and Φ_{f} is a linear matrix of BD.

Similarly, we will also use the recursive least square estimation method to calculate *BD* recursively. Let $\boldsymbol{z}_{\mathbf{f}} = [\boldsymbol{z}_{1}^{\mathrm{T}} \ \boldsymbol{z}_{2}^{\mathrm{T}} \cdots \boldsymbol{z}_{f}^{\mathrm{T}}]^{\mathrm{T}}$ be a vector as the form of $\boldsymbol{u}_{f}(k)$, namely, $\boldsymbol{z}_{j} \in \mathbf{R}^{m}$, $j = 1, 2, \cdots, f$.

Let

$$\boldsymbol{\theta} = \mathbf{Vec}(BD)$$

where BD is the product of B and D. Denote

$$M(C, A, \boldsymbol{z}_{f}) = \begin{bmatrix} \boldsymbol{0} \\ \boldsymbol{z}_{1}^{\mathrm{T}} \otimes C \\ \boldsymbol{z}_{1}^{\mathrm{T}} \otimes CA + \boldsymbol{z}_{2}^{\mathrm{T}} \otimes C \\ \vdots \\ \sum_{i=1}^{f-1} \boldsymbol{z}_{i}^{\mathrm{T}} \otimes CA^{f-1-i} \end{bmatrix}$$

where \otimes is Kronecker product. For any vector $\boldsymbol{u}_{f}(k)$ as defined in Section 1,

$$\Phi_f \boldsymbol{u}_f(k) = M(C, A, \boldsymbol{u}_f(k))\boldsymbol{\theta} + \boldsymbol{\omega}_k$$
(19)

where $\boldsymbol{\omega}_k$ is the estimation noise. Denote

$$\widetilde{O}_{f}(k) = \begin{bmatrix} \widehat{C}(k) \\ \widehat{C}(k)\widehat{A}(k) \\ \vdots \\ \widehat{C}(k)\widehat{A}^{f-1}(k) \end{bmatrix}$$
$$\widehat{O}_{f}^{\perp}(k) = (\widetilde{O}_{f}(k))^{\perp}$$
$$M(k) = M(\widehat{C}(k), \widehat{A}(k), u_{f}(k))$$
$$\psi(k) = \left[(\widehat{O}_{f}^{\perp}(k))^{\mathrm{T}}M(k) \right]^{\mathrm{T}} = M^{\mathrm{T}}(k)\widehat{O}_{f}^{\perp}(k)$$

where $O_f^{\perp}(k)$ is the orthogonal complement matrix of k-th extended observable estimated matrix $\tilde{O}_f(k)$ generated by $\hat{C}(k)$ and $\hat{A}(k)$, i.e., satisfying $\widetilde{O}_f(k)^{\mathrm{T}} O_f^{\perp}(k) = 0$, and $\widehat{O_{f}^{\perp}}(k) \in \mathbf{R}^{lf \times lf - n}$ is full column rank. M(k) is the k-th matrix substituting $\hat{C}(k)$ and $\hat{A}(k)$ into $M(C, A, \boldsymbol{u}_f(k))$ instead of C and A.

Next, we put (19) into (18) and gain the following result:

$$(O_f^{\perp})^{\mathrm{T}} \boldsymbol{y}_f(k+1) = (O_f^{\perp})^{\mathrm{T}} M(C, A, \boldsymbol{u}_f(k)) \boldsymbol{\theta} + (O_f^{\perp})^{\mathrm{T}} \boldsymbol{e}_f(k)$$
(20)

Then $\boldsymbol{\theta}$, namely, **Vec**(*BD*) can be estimated using recursive least square method as follows:

$$\begin{cases} \psi(k) = (O_f^{\perp})^{\mathrm{T}} M(C, A, \boldsymbol{u}_f(k)) \\ \hat{\boldsymbol{\theta}}(k+1) = \boldsymbol{\theta}(k) + \kappa(k) P(k) \psi(k) [(O_f^{\perp})^{\mathrm{T}} \boldsymbol{y}_f(k+1) - (O_f^{\perp})^{\mathrm{T}} M(C, A, \boldsymbol{u}_f(k)) \boldsymbol{\theta}] \\ P(k+1) = P(k) + \kappa(k) P(k) \psi(k) \psi(k)^{\mathrm{T}} P(k) \\ \text{here} \end{cases}$$

where

$$P(k) = \left(\sum_{i=0}^{k-1} \psi(k)\psi(k)^{\mathrm{T}}\right)^{-1} \\ \kappa(k) = (1 + \psi(k)^{\mathrm{T}}G(k)\psi(k))^{-1}$$

The initial values P(0), G(0) could be chosen as I with proper dimensions, and $A(0), \boldsymbol{\theta}(0)$ could be taken as any deterministic matrix of compatible dimensions.

Once $\operatorname{Vec}(BD)$ is obtained, since $D_1 = I$ by the definition of D we could observe B directly from BD, i.e., B = BD(1 : n, 1 : m). And then by the property of the Kronecker product $ABC = C^{\mathrm{T}} \otimes B^{\mathrm{T}} \otimes A^{\mathrm{T}}$ with arbitrary suitable dimension, respectively, D could be estimated uniquely under proper restriction on B.

3 Simulation results

In this section, we will give a simulated example for the case of MIMO Hammerstein system to evaluate the proposed algorithm in the previous section:

$$\boldsymbol{x}_{k+1} = A\boldsymbol{x}_k + Bh(\boldsymbol{u}_k) \tag{21}$$

$$\boldsymbol{y}_{k+1} = C \boldsymbol{x}_{k+1} + \boldsymbol{e}_{k+1}$$
 (22)

where

$$A = \begin{bmatrix} 0 & 1 & 0 \\ -0.5625 & -1.5 & 0 \\ 0 & 0 & 0.85 \end{bmatrix}$$
$$B = \begin{bmatrix} 1.5 & 0.5 \\ 0 & 1.5 \\ 0.1 & 0.2 \end{bmatrix}$$
$$C = \begin{bmatrix} 0 & 0.3 & 0 \\ 0.5 & 0 & 2 \end{bmatrix}$$
$$D_1 = I_{2 \times 2}$$
$$D_2 = I_{2 \times 2}$$

The static nonlinearity is chosen as

$$\boldsymbol{h}(u) = \begin{bmatrix} 1 + \sin u(1) \\ \cos\left(\frac{u(2)}{2}\right) + 2\sin\left(\frac{u(2)}{2}\right) \end{bmatrix}$$

and the inputs u(i), i = 1, 2 are taken uniformly on the interval [-1, 1]. The observed noise is also chosen as a sequence of Gaussian white noise with standard deviation $\sigma = 0.1 \times I$, where data set has 3 000 input and output measurements respectively. Let $N_0 = 1000$.

To evaluate the gap between the spanned subspace by the column vector of \hat{O}_f and O_f , we introduce the following variant:

$$\Theta(\operatorname{span}(O_f), \operatorname{span}(\hat{O}_f(k))) = \|\Pi_{O_f} - \Pi_{\hat{O}_f(k)}\|$$

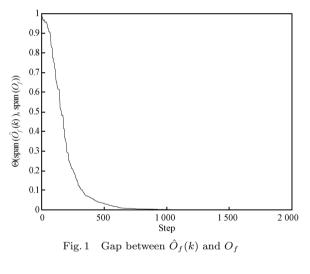
where \prod_X denotes the orthogonal projection matrix on the spanned space of column of X. By the definition of Θ , it can be considered as an angle between the two subspaces. The smaller the angle is, the closer the two subspaces will be.

As shown in Fig.1, the gap between the true value O_f and its estimation $\hat{O}_f(k)$ has approached to zero as soon as it turns to the fast averaging computation on the 1 000-th step.

Fig. 2 has given a comparison in the end of steps from 1500 to 2000 between the normal PCA method in the winding line and method using fast averaging technique underneath in the straight line. We can observe that the line in the fast averaging method goes to zero more smoothly and quickly than the normal method. It is indicated that fast averaging technique is a more satisfying choice in identifying the extended observable matrix.

The true eigenvalues of A are 0.85, -0.75, and -0.75. Fig. 3 has shown that the estimation of A's eigenvalues converge to its true value gradually as the computation process goes to the end.

Fig. 4 has indicated that the least square estimation of B and D could be rather accurate.



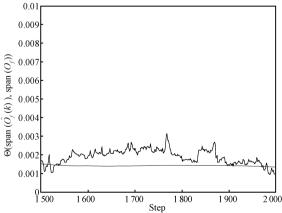


Fig. 2 Comparison between fast averaging method and normal method

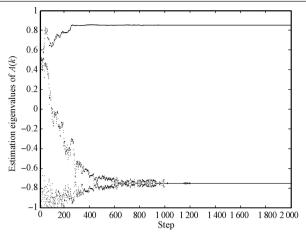


Fig. 3 Estimation eigenvalues of A(k)

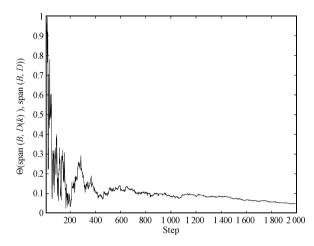


Fig. 4 Gap between the two subspaces spanned by B, D and their true values

4 Conclusion

In this paper, a method for estimating the parameters of MIMO Hammerstein system in state-space form applying subspace identification method was proposed. Extension of Hammerstein system to the MIMO case has brought in many difficulties. The existing methods mainly deal with SISO case, among which the relation between the input and output is indicated in the transfer function form. However, the transfer function in the MIMO form is difficult to show and the relationship between the parameters and the inputs is troublesome to indicate. So, the state space model is applied. In particular, the presented recursive algorithm based on PCA in view of the stochastic approximation methods has been given. This is an important extension of the previous works, which have been restricted to the area of SISO Hammerstein system. With a purpose to speed up the computation rate, we have applied the PCA algorithm based on the stochastic approximation using fast averaging technique.

Although the results obtained here are satisfying, to extend the subspace identification methods on the general nonlinear system, further research need to be performed. To recursively identify this general nonlinear system is also our future focus of research.

Appendix A The proof of Proposition 1

To show $r_{xx} - r_{xu}R_{uu}^{-1}r_{ux} > 0$, it is equivalent to prove that

$$\mathbf{E}\left\{\begin{bmatrix}\boldsymbol{x}_{k}\\\boldsymbol{u}_{f}(k)\end{bmatrix}\times\begin{bmatrix}\boldsymbol{x}_{k}\\\boldsymbol{u}_{f}(k)\end{bmatrix}^{\mathrm{T}}\right\} = \begin{bmatrix}r_{xx} & r_{xu}\\r_{uu} & r_{uu}\end{bmatrix} > 0 \qquad (A1)$$

By the state space model (7) and (8). Let $\boldsymbol{x}_k = \frac{F^{\boldsymbol{u}}(q^{-1})}{a(q^{-1})} \tilde{\boldsymbol{u}}_k$, where q^{-1} is the backward shift operator, i.e., $q^{-1} \tilde{\boldsymbol{u}}_k = \tilde{\boldsymbol{u}}_{k-1}$,

$$F^{\mathbf{u}}(q^{-1}) = \frac{\operatorname{Adj}(qI - A)BD}{q^n} = F_1^{\mathbf{u}}q^{-1} + F_2^{\mathbf{u}}q^{-2} + \dots + F_n^{\mathbf{u}}q^{-n}$$
$$a(q^{-1}) = \frac{\det(qI - A)}{q^n} = a_0 + a_1q^{-1} + a_2q^{-2} + \dots + a_nq^{-n}$$

where $F_i^{\boldsymbol{u}} \in \mathbf{R}^{n \times Nm}, a_0 = 1.$

$$\begin{bmatrix} \boldsymbol{x}_k \\ \boldsymbol{u}_f(k) \end{bmatrix} = \frac{\Xi}{a(q^{-1})} \begin{bmatrix} \tilde{\boldsymbol{u}}_{k-n} \\ \vdots \\ \tilde{\boldsymbol{u}}_{k+f-1} \end{bmatrix}$$

where

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$$\Xi = \begin{bmatrix} F_n^{\mathbf{u}} & F_{n-1}^{\mathbf{u}} & \cdots & F_1^{\mathbf{u}} & 0 & 0 & \cdots & 0\\ 0 & 0 & \cdots & 0 & a_n I & a_{n-1} I & \cdots & a_0 I\\ 0 & 0 & \cdots & a_n I & a_{n-1} I & a_{n-2} I & \cdots & 0\\ \vdots & \vdots & \ddots & \vdots & \vdots & \vdots & \ddots & \vdots\\ 0 & a_n I & \cdots & a_2 I & a_1 I & a_0 I & \cdots & 0\\ a_n I & a_{n-1} I & \cdots & a_1 I & a_0 I & 0 & \cdots & 0 \end{bmatrix}$$

If Ξ has full row rank under A 1 ~ A 4, and

$$\mathbf{E} \left\{ \frac{1}{a(q^{-1})} \begin{bmatrix} \tilde{\boldsymbol{u}}_{k-n} \\ \vdots \\ \tilde{\boldsymbol{u}}_{k+f-1} \end{bmatrix} \times \frac{1}{a(q^{-1})} \begin{bmatrix} \tilde{\boldsymbol{u}}_{k-n} \\ \vdots \\ \tilde{\boldsymbol{u}}_{k+f-1} \end{bmatrix}^{\mathrm{T}} \right\}$$
(A2)

is positive definite, then (A1) is true. Note that

$$\begin{bmatrix} \tilde{\boldsymbol{u}}_{k-n} \\ \vdots \\ \tilde{\boldsymbol{u}}_{k+f-1} \end{bmatrix} = \begin{bmatrix} \boldsymbol{\varphi}^{\mathrm{T}}(\boldsymbol{u}_{k-n}) \\ \vdots \\ \boldsymbol{\varphi}^{\mathrm{T}}(\boldsymbol{u}_{k+f-1}) \end{bmatrix}$$

is persistently excited as stated in the remark after the assumptions A 1 \sim A 4.

At last, we will show that Ξ is full row rank. Equivalently, we can show that $[F_n^u \cdots F_1^u]$ is full rank. Define

$$F^{u}(q^{-1}) = \frac{\operatorname{Adj}(qI - A)BD}{q^{n}} = a(q^{-1})(qI - A)^{-1}BD = (a_{0} + a_{1}q^{-1} + \dots + a_{n}q^{-n})q^{-1}$$
$$(I + q^{-1}A + q^{-2}A^{2} + \dots)BD$$

1

Comparing the coefficients of the two sides of this equation,

we have

$$[F_n^{u} \cdots F_1^{u}] = [A^{n-1}BD \ A^{n-2}BD \cdots BD] \times \begin{bmatrix} a_0I & 0 & \cdots & 0\\ a_1I & a_0I & \cdots & 0\\ \vdots & \vdots & \ddots & \vdots\\ a_{n-1}I & a_{n-2}I & \cdots & a_0I \end{bmatrix}$$
(A3)

Since $[A^{n-1}BD, A^{n-2}BD, \cdots, BD]$ is full rank in column by the controllability of (A; B) in assumption A 1, $D = [I, D_2, \dots, D_N]$, and the right matrix consisting of a_0, a_1, \dots, a_N of the right equality of (A3) is nonsingular, $[F_n^{\boldsymbol{u}}, \dots, F_1^{\boldsymbol{u}}]$ is also full rank. \Box

Appendix B The proof of Theorem 1

From the analysis in Subsection 1.3, we know $\operatorname{span}\{O_f\} = \operatorname{span}\{R_s\}$, by Theorem 1, under the assumptions A 1 ~ A 5, then we will show $\hat{O}_f(k) \xrightarrow[k \to \infty]{a.s.} O_f T$, under the assumptions A 1 \sim A 5.

Since $\{\boldsymbol{u}_k\}$ is an independent identically distributed process, $\{\boldsymbol{\varphi}_i, i = 1, \cdots, N\}$ are bounded, $\{\tilde{\boldsymbol{u}}_k\}$ is a bounded process, the observed noise $\{e_k\}$ is bounded, Y_k is bounded. By the definition of $\hat{\boldsymbol{z}}_f(k)$, then $R_k = \mathbb{E}\{\hat{\boldsymbol{z}}_f(k)\hat{\boldsymbol{z}}_f(k)^{\mathrm{T}}\}$ is also bounded. In view of Chapter 3.3 of [15], we just have to show the noise condition:

$$\lim_{T \to 0} \limsup_{k \to \infty} \frac{1}{T} \left\| \sum_{i=k}^{m(k,T)} a_k (\hat{\boldsymbol{z}}_f(i+1) \hat{\boldsymbol{z}}_f^{\mathrm{T}}(i+1) - R) \right\| = 0,$$
$$\forall T_k \in (0,T]$$

where $m(k,T) = \max\{m, \sum_{i=k}^{m} \leq T\}$. Without loss of generality, we assume $a_k = 1/k$. By the result of Lemma 1,

$$\frac{1}{k} (U_k U_k^{\mathrm{T}})^{-1} \xrightarrow[k \to \infty]{\text{a.s.}} R_{uu}^{-1}$$
$$\frac{1}{k} (Y_k U_k^{\mathrm{T}}) \xrightarrow[k \to \infty]{\text{a.s.}} R_{yu}$$

Then, we use the second equation to subtract the first one and it turns out to be:

$$\Delta(k) = R_{yu}R_{uu}^{-1} - \left(\frac{1}{k}Y_kU_k^{\mathrm{T}}\right)\frac{1}{k}(U_kU_k^{\mathrm{T}})^{-1} \xrightarrow[k \to \infty]{a.s.} 0$$

$$\begin{aligned} \hat{\boldsymbol{z}}_{f}(i+1)\hat{\boldsymbol{z}}_{f}^{T}(i+1) - R &= \\ [\boldsymbol{y}_{f}(i+1) - Y_{i+1}U_{i+1}^{T}(U_{k}U_{k}^{T})^{-1}\boldsymbol{u}_{f}(i+1)] \times \\ [\boldsymbol{y}_{f}(i+1) - Y_{i+1}U_{i+1}^{T}(U_{k}U_{k}^{T})^{-1}\boldsymbol{u}_{f}(i+1)]^{T} - R &= \\ [\boldsymbol{y}_{f}(i+1) - Y_{i+1}U_{i+1}^{T}(U_{k}U_{k}^{T})^{-1}\boldsymbol{u}_{f}(i+1) + \Delta(k)\boldsymbol{u}_{f}(i+1)] \times \\ [\boldsymbol{y}_{f}(i+1) - Y_{i+1}U_{i+1}^{T}(U_{k}U_{k}^{T})^{-1}\boldsymbol{u}_{f}(i+1)]^{T} + \\ \Delta(k)\boldsymbol{u}_{f}(i+1)^{T} - R &= \\ [\hat{\boldsymbol{z}}_{f}(i+1) + \Delta(i+1)\boldsymbol{u}_{f}(i+1)]^{T} - R &= \\ [\hat{\boldsymbol{z}}_{f}(i+1) + \Delta(i+1)\boldsymbol{u}_{f}(i+1)]^{T} - R &= \\ \hat{\boldsymbol{z}}_{f}(i+1)\hat{\boldsymbol{z}}_{f}(i+1)^{T} - R + \hat{\boldsymbol{z}}_{f}(i+1)\boldsymbol{u}_{f}^{T}(i+1)\Delta^{T}(i+1) + \\ \Delta(i+1)\boldsymbol{u}_{f}(i+1)\boldsymbol{z}_{f}^{T}(i+1) + \Delta^{T}(i+1)\boldsymbol{u}_{f}^{T}(i+1)\Delta^{T}(i+1) \end{aligned}$$

As we have known that $\Delta(i+1) \xrightarrow[i \to \infty]{\text{a.s.}} 0$ and $[\boldsymbol{u}_f(i+1)]$

1), $\hat{\boldsymbol{z}}_f(i+1)$] is bounded, then

$$\begin{split} \lim_{T \to 0} \limsup_{k \to \infty} \frac{1}{T} \bigg\| \sum_{i=k}^{m(k,T)} \frac{1}{i} [\hat{\boldsymbol{z}}_f(i+1) \boldsymbol{u}_f^{\mathrm{T}}(i+1) \Delta^{\mathrm{T}}(i+1) + \\ \Delta(i+1) \boldsymbol{u}_f(i+1) \boldsymbol{z}_f^{\mathrm{T}}(i+1) + \Delta^{\mathrm{T}}(i+1) \boldsymbol{u}_f^{\mathrm{T}}(i+1) \Delta^{\mathrm{T}}(i+1)] \bigg\| = 0 \\ \text{a.s. } \forall \ T_k \ \in (0,T] \end{split}$$

and we just have to show that

$$\lim_{T \to 0} \limsup_{k \to \infty} \frac{1}{T} \parallel \sum_{i=k}^{m(k,T)} \frac{1}{i} [\hat{\boldsymbol{z}}_f(i+1) \hat{\boldsymbol{z}}_f^{\mathrm{T}}(i+1) - R] \parallel = 0,$$

a.s. $\forall T_k \in (0,T]$

From Lemma 1, we have come to know that $\{(\boldsymbol{y}_{k}^{\mathrm{T}}, \boldsymbol{u}_{k}^{\mathrm{T}})^{\mathrm{T}}\}$ has the stable ergodic characteristic, so does $\{\boldsymbol{z}_{f}(k)\boldsymbol{z}_{f}(k)^{\mathrm{T}}\}$. By (11), $\mathrm{E}\{\tilde{R}(k)\} = \mathrm{E}\{\hat{\boldsymbol{z}}_{f}(k)\hat{\boldsymbol{z}}_{f}(k)^{\mathrm{T}}\} = R$ and $\frac{1}{k}\sum_{i=1}^{k}\tilde{R}(k) - R \xrightarrow[i \to \infty]{i \to \infty} 0$. Let $\eta_{k} = \sum_{i=1}^{k}(\tilde{R}(k) - R)$, i.e., $\eta_{k} \xrightarrow[i \to \infty]{i \to \infty} 0$,

$$\begin{split} &\sum_{i=k}^{m(k,T)} \frac{1}{i} (\hat{\mathbf{z}}_f(i+1) \hat{\mathbf{z}}_f^{\mathrm{T}}(i+1) - R) = \\ &\sum_{i=k}^{m(k,T)} \frac{1}{i} (\eta_{i+1} - \eta_i) = \\ &\frac{1}{m(k,T)} \eta_{m(k,T)+1} - \frac{1}{k} \eta_k + \sum_{i=k+1}^{m(k,T)} \eta_k (\frac{1}{i+1} - \frac{1}{i}) = \\ &\frac{1}{m(k,T)} \eta_{m(k,T)+1} - \frac{1}{k} \eta_k + \sum_{i=k+1}^{m(k,T)} \eta_k (\frac{1}{(i+1)i}) \\ &\xrightarrow[i \to \infty]{} 0 \end{split}$$

So, (4) holds naturally.

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