Conditioning analysis of a continuous time subspace-based model identification algorithm

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Abstract

We present in this paper a study concerning the conditioning analysis of a continuoustime deterministic subspace-based model identification algorithm. We show that the conditioning number of the associated extended observability matrices depends on an exponential way from both: the estimated order of the system and the dimension of the system output vector.

1 Introduction

As far as linear time-invariant multivariable systems are concerned, subspace-based model identification algorithms constitute a broad family of identification methods mainly characterized by the use of geometric information (see for instance [3], [4], and [5]). The subspace-based methods compute the estimated parameters, i.e., a state space realization of the system, from an approximation of the observability subspace of the concerned system. This approximation is obtained from a discrete-time set of input-output measurements. Continuous-time data is filtered and sampled in order to obtain discrete-time information (see [1]). We present in this paper an algorithm to identify a continuous linear time-invariant model from a given set of discrete-time input-output measurements. Our algorithm is based on the method proposed in [3], and depends essentially on the pseudo-inverse of the so-called extended observability matrices. We present then the analysis of the conditioning number of these observability matrices, which is the main purpose of this paper.

In Section 2 we present our proposed algorithm, while in Section 3 we discuss its numerical properties.

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2 The CN4SID algorithm

Consider the continuous linear time-invariant system (A, B, C, D) described by:

$$\begin{cases} \mathbf{p}x(t) = Ax(t) + Bu(t), \\ y(t) = Cx(t) + Du(t), \end{cases}$$
(2.1)

where: \mathbf{p} denotes the differential operator, i.e. $\mathbf{p} = d/dt$; $x(\cdot) \in \mathbb{R}^n$ denotes the state; $u(\cdot) \in \mathbb{R}^m$ denotes the input, and $y(\cdot) \in \mathbb{R}^p$ denotes the output. A, B, C and D are linear maps represented by real constant matrices. It is assumed that (C, A) is observable.

Consider the scalar causal stable operator:

$$\lambda = \frac{1}{1 + \mathbf{p}\tau},\tag{2.2}$$

where τ is a scalar such that $\tau > 0$. Given a \mathbb{R}^p valued signal y(t), we define the i^{th} -filtered signal $[\lambda^i y](t)$ as follows:

$$\left[\lambda^{i}y\right]\left(t\right) = \begin{cases} \lambda \left[\lambda^{i-1}y\right]\left(t\right), & \text{for } i > 1, \\ y\left(t\right), & \text{when } i = 0. \end{cases}$$

Applying this filtering action to the system (A, B, C, D) we obtain the modified system:

$$\begin{cases} x(t) = A_{\lambda} [\lambda x](t) + B_{\lambda} [\lambda u](t), \\ y(t) = Cx(t) + Du(t), \end{cases}$$
(2.3)

where $A_{\lambda} = I + \tau A$ and $B_{\lambda} = \tau B$. Since (A, B) and $(A_{\lambda}, B_{\lambda})$ are linked through a biyective map, the obtention of estimates for A and B results in the obtention of estimates for A and B.

Suppose constant sampling and a sampling time sequence $\{t_k\}_{k=0}^{N-1}$ be given, where $t_k = t_0 + kh$ (with h denoting the sampling period). Let input-output measurements $\{u_k\}_{k=0}^{N-1}$ and $\{y_k\}_{k=0}^{N-1}$ on system (2.3) be given. The deterministic Continuous Subspace Identification (CSId) problem is then defined as follows:

Definition 2.1. CSId problem: Consider the filtered model (2.3) and the input-out measurements $\{u_k\}_{k=0}^{N-1}$ and $\{y_k\}_{k=0}^{N-1}$, estimate then the unknown matrices A, B, C, and D.

We can then build the following matrices:

$$Y := \begin{bmatrix} [\lambda^{i-1}y]_0 & [\lambda^{i-1}y]_1 & \cdots & [\lambda^{i-1}y]_{N-1} \\ [\lambda^{i-2}y]_0 & [\lambda^{i-2}y]_1 & \cdots & [\lambda^{i-2}y]_{N-1} \\ \vdots & \vdots & \ddots & \vdots \\ y_0 & y_1 & \cdots & y_{N-1} \end{bmatrix}$$

$$(2.4)$$

and:

$$U := \begin{bmatrix} [\lambda^{i-1}u]_0 & [\lambda^{i-1}u]_1 & \cdots & [\lambda^{i-1}u]_{N-1} \\ [\lambda^{i-2}u]_0 & [\lambda^{i-2}u]_1 & \cdots & [\lambda^{i-2}u]_{N-1} \\ \vdots & \vdots & \ddots & \vdots \\ u_0 & u_1 & \cdots & u_{N-1} \end{bmatrix}.$$

$$(2.5)$$

In what follows we present the **CN4SID** algoritm, which is based on: the **N4SID** algorithm [3], and the space transformation discussed in [1].

CN4SID Algorithm: Consider a continous linear-time invariant system described by (2.1) and the corresponding filtered model (2.3). Let the input-output sequences $\{u_k\}_{k=0}^{N-1}$ and $\{y_k\}_{k=0}^{N-1}$ be given. Let i be an the estimated order of the system such that N >> i > n, where n denotes the real system. It is then possible to find matrices \widehat{A} , \widehat{B} , \widehat{C} , and \widehat{D} (i.e., the estimates of A, B, C, and D, respectively), according to the procedure:

- 1. Build matrices Y and U as defined in (2.4) and (2.5), respectively.
- 2. Build matrices $[\lambda^i Y]$ and $[\lambda^i U]$.
- 3. Build matrices $[\lambda^i W]$ and $[\lambda^i W]$ defined as follows:

$$\left[\lambda^{i}W\right]:=\left[\begin{array}{c}\left[\lambda^{i}U\right]\\\left[\lambda^{i}Y\right]\end{array}\right],\ \mathrm{and}\ \left[\lambda^{i+1}W\right]:=\left[\begin{array}{c}\left[\lambda^{i+1}U\right]\\\left[\lambda^{i+1}Y\right]\end{array}\right].$$

4. Compute the proyections:

$$Y/\left[\begin{array}{c} [\lambda^{i}W] \\ [\lambda^{i}U] \end{array}\right] = \left[\begin{array}{c} L_{w} & L_{u} \end{array}\right] \left[\begin{array}{c} [\lambda^{i}W] \\ [\lambda^{i}U] \end{array}\right]$$

and:

$$\left[\lambda Y\right] / \left[\begin{array}{c} \left[\lambda^{i+1} W\right] \\ \left[\lambda^{i+1} U\right] \end{array} \right] = \left[\begin{array}{c} L_{\overline{w}} \end{array} \right] \left[\begin{array}{c} \left[\lambda^{i+1} W\right] \\ \left[\lambda^{i+1} U\right] \end{array} \right].$$

5. From the singular value decomposition of L_w :

$$\Lambda_i X_i = L_w \begin{bmatrix} \lambda^i W \end{bmatrix} = \begin{bmatrix} U_N & U_0 \end{bmatrix} \begin{bmatrix} S_n & 0 \\ 0 & S_0 \end{bmatrix} \begin{bmatrix} V_n' \\ V_0' \end{bmatrix} \begin{bmatrix} \lambda^i W \end{bmatrix}$$

compute the i^{th} - Λ -extended observability matrices:

$$\Lambda_i = U_n S_n^{1/2},\tag{2.6}$$

where $S_n^{1/2}$ stands for the Cholesky factor of S_n .

6. From extended observability matrices Λ_i compute the state sequences:

$$X_i = \Lambda_i^{\dagger} L_w \left[\lambda^i W \right]$$

and:

$$[\lambda X_i] = \Lambda_i^{\dagger} L_{\overline{w}} \left[\lambda^{i+1} W \right].$$

 Λ_i^{\dagger} stands for the pseudo-inverse of Λ_i .

7. Solve the following optimization problem:

$$\min_{k} \left\| \left[egin{array}{c} X_i \ [\lambda^i y] \end{array}
ight] ext{-} K \left[egin{array}{c} [\lambda X_i] \ [\lambda^i u] \end{array}
ight]
ight\|_F^2,$$

where:

$$K := \left[\begin{array}{cc} A_{\lambda} & B_{\lambda} \\ \widehat{C} & \widehat{D} \end{array} \right].$$

and $\|\alpha\|_F$ stands for the Frobenius norm of α .

8. Finally:

$$\widehat{A} = \frac{A_{\lambda} - I}{\tau}$$

and:

$$\widehat{B} = \frac{B_{\lambda}}{\tau}.$$

In what follows we discuss the numerical properties of the CN4SID algoritm.

3 Numerical properties

Since the **CN4SID** algoritm depends essentially on the pseudo-inverse of the Λ -extended observability matrices Λ_i given by (2.6), in this section we discuss the conditioning of these matrices.

We first present some preliminary results concerning the nature of Λ_i . Let us define the i^{th} -extended observability matrix Γ_i as follows:

$$\Gamma_{i} := \begin{bmatrix} C \\ CA \\ CA^{2} \\ \vdots \\ CA^{i\cdot 1} \end{bmatrix}. \tag{3.7}$$

Then:

Lemma 3.1. Let Γ_i be as defined in (3.7), then:

$$\Lambda_i = T\Gamma_i$$

where:

$$T = \begin{bmatrix} I & 0 & 0 & \cdots & 0 \\ I & \tau I & 0 & \cdots & 0 \\ I & 2\tau I & \tau^{2} I & \cdots & 0 \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ I & \binom{i-1}{1}\tau I & \binom{i-1}{2}\tau^{2} I & \cdots & \tau^{i-1} I \end{bmatrix} \in \mathbb{R}^{ip \times ip}, \tag{3.8}$$

with:

$$\binom{v}{z} := \frac{v!}{(v-z)!z!}.$$

Proof. Because of the definition of Λ_i we have that:

$$\Lambda_{i} = \begin{bmatrix} C \\ CA_{\lambda} \\ \vdots \\ CA_{\lambda}^{i-1} \end{bmatrix},$$
(3.9)

where:

$$A_{\lambda} = I + \tau A. \tag{3.10}$$

Thus:

$$\Lambda_{i} = \begin{bmatrix}
C \\
C + \tau C A \\
C + 2\tau C A + \tau^{2} C A^{2} \\
\vdots \\
C + \binom{i \cdot 1}{1} \tau C A + \binom{i \cdot 1}{2} \tau^{2} C A^{2} + \dots + \tau^{i \cdot 1} C A^{i \cdot 1}
\end{bmatrix}$$
(3.11)

which implies:

$$\Lambda_i = T\Gamma_i. \tag{3.12}$$

Thus, the numerical conditioning of Λ_i , in terms of the Frobenius norm, is given by:

$$k_{F}\left(\Lambda_{i}
ight)=\left\Vert \Lambda_{i}\right\Vert _{F}\left\Vert \Lambda_{i}^{ ext{-}1}\right\Vert _{F},$$

and from Lemma 3.1:

$$k_F(\Lambda_i) = \|T\Gamma_i\|_F \|(T\Gamma_i)^{-1}\|_F.$$

Consequently:

$$k_F(\Lambda_i) \leq ||T||_F ||T^{-1}||_F ||\Gamma_i||_F ||\Gamma_i^{-1}||_F,$$

i.e.:

$$k_F(\Lambda_i) \le k_F(T) k_F(\Gamma_i). \tag{3.13}$$

Remark 3.1. Because of its dependence on the input-output measurements, it is not possible to have a priori knowledge of $k_F(\Gamma_i)$. However, it is suitable to have a small value for this conditioning number. As far as $k_F(T)$ is concerned, it is possible to have a priori knowledge of its value. For the above reason, we focus our analysis on the study of the conditioning number of T.

Lemma 3.2. Let T be as defined in (3.8), then:

$$||T||_F^2 = p \sum_{k=1}^i \sum_{j=1}^k {k-1 \choose j-1}^2 \tau^{2(j-1)}.$$

Proof. Let matrix T be as defined in (3.8). Then, the block matrix $T_{kj} \in \mathbb{R}^{p \times p}$ can be written as follows:

$$T_{kj} = {\binom{k-1}{j-1}} \tau^{j-1} I_{p \times p} \tag{3.14}$$

and its corresponding transpose block matrix is given by:

$$T'_{jl} = {l-1 \choose j-1} \tau^{j-1} I_{p \times p} \tag{3.15}$$

and consequently the kl-block matrix of TT', i.e. $(TT')_{kl} \in \mathbb{R}^{p \times p}$, is given by:

$$(TT')_{kl} = \sum_{j=1}^{i} {k-1 \choose j-1} \tau^{j-1} I_{p \times p} {l-1 \choose j-1} \tau^{j-1} I_{p \times p}$$

$$= \sum_{j=1}^{i} {k-1 \choose j-1} {l-1 \choose j-1} \tau^{2(j-1)} I_{p \times p}.$$
(3.16)

Moreover:

$$tr(TT') = \sum_{k=l=1}^{i} \sum_{j=1}^{i} {k-1 \choose j-1} {l-1 \choose j-1} \tau^{2(j-1)} p$$

$$= \sum_{k=1}^{i} \sum_{j=1}^{k} {k-1 \choose j-1}^{2} \tau^{2(j-1)} p,$$
(3.17)

and because of the definition of the Frobenius norm:

$$||T||_F^2 = tr(TT') \tag{3.18}$$

$$= p \sum_{k=1}^{i} \sum_{j=1}^{k} {\binom{k-1}{j-1}}^2 \tau^{2(j-1)}$$
(3.19)

Remark 3.2. From the previous lemma we have that $||T||_F^2$ is a 2(j-1) polynomial function of τ . It is usual that $0 < \tau < 1$. Thus, $||T||_F^2$ usually has a small value.

As far as the Frobenius norm of T^{-1} is concerned, we have the following result, which gives the Frobenius norm of adj(T), i.e., the adjunct of T:

Lemma 3.3. Let T be as defined in (3.8), then:

$$\|adj(T)\|_F^2 = p\tau^{n(n-1)} \sum_{k=1}^i \sum_{j=1}^k {k-1 \choose j-1}^2 \tau^{-2(k-1)},$$

where p > 0 denotes the dimension of the output vector.

Proof. Let matrix T be as defined in (3.8). Then, the adjunct of T can be written in terms of its kj-block matrix, i.e. $(adjT)_{kj} \in \mathbb{R}^{p \times p}$, as follows:

$$(adjT)_{kj} = (-1)^{k+j} {\binom{k-1}{j-1}} \tau^{\frac{i(i-1)}{2} \cdot (k-1)} I_{p \times p}$$
(3.20)

and the corresponding block matrices of the transpose of adj(T) are given by:

$$(adjT)'_{jl} = (-1)^{l+j} {l-1 \choose j-1} \tau^{\frac{i(i-1)}{2} - (l-1)} I_{p \times p}.$$
(3.21)

Then, the block matrices of the matrix product adj(T) adj(T)' can be written as:

$$(adjT (adjT'))_{kl} = \sum_{j=1}^{i} (-1)^{k+j} {k-1 \choose j-1} \tau^{\frac{i(i-1)}{2} - (k-1)} I_{p \times p}$$

$$\cdot (-1)^{l+j} {l-1 \choose j-1} \tau^{\frac{i(i-1)}{2} - (l-1)} I_{p \times p}$$

$$= \sum_{j=1}^{i} (-1)^{k+l+2j} {k-1 \choose j-1} {l-1 \choose j-1} \tau^{i(i-1)-2(k-1)} I_{p \times p}.$$

$$(3.22)$$

As far as the trace of the matrix product adj(T) adj(T)' is concerned we have:

$$tr\left(adj\left(T\right)adj\left(T\right)'\right) = \sum_{k=l=1}^{i} \sum_{j=1}^{i} (-1)^{k+l+2j} \binom{k-1}{j-1}$$

$$\cdot \binom{l-1}{j-1} \tau^{\frac{i(i-1)}{2} - (k-1)} \tau^{\frac{i(i-1)}{2} - (l-1)} p$$

$$= \sum_{k=1}^{i} \sum_{j=1}^{i} \binom{k-1}{j-1}^{2} \tau^{i(i-1)-2(k-1)} p. \tag{3.23}$$

Consequently:

$$||adj(T)||_{F}^{2} : = tr \left(adj(T) adj(T)'\right)$$

$$= p\tau^{i(i-1)} \sum_{k=1}^{i} \sum_{j=1}^{k} {k-1 \choose j-1}^{2} \tau^{-2(k-1)}.$$
(3.24)

We can at this level present our main result. Combining Lemma 3.2 and Lemma 3.3 we have:

Theorem 3.1. Let T be as defined in (3.8), then:

$$k_{F}\left(T\right) = p\tau^{\frac{i(i-1)(1-p)}{2}}\sqrt{\left(\sum_{k=1}^{i}\sum_{j=1}^{k}\binom{k-1}{j-1}^{2}\tau^{2(j-1)}\right)\left(\sum_{k=1}^{i}\sum_{j=1}^{k}\binom{k-1}{j-1}^{2}\tau^{-2(j-1)}\right)}.$$

Proof. By definition:

$$k_F(T) := \|T\|_F \|T^{-1}\|_F,$$
 (3.25)

i.e.:

$$k_{F}(T) = \sqrt{tr(TT')tr(T^{-1}(T^{-1})')}$$

$$= \sqrt{tr(TT')\frac{tr(adj(T)adj(T)')}{\det(TT')}}.$$
(3.26)

Now, taking into account Lemma 3.2, Lemma 3.3, and since $\det(T) = \tau^{\frac{ip(i-1)}{2}}$, we have:

$$k_F(T)$$

$$= \sqrt{\left(p\sum_{k=1}^{i}\sum_{j=1}^{k} {k-1 \choose j-1}^{2} \tau^{2(j-1)}\right) \left(p\tau^{i(i-1)}\sum_{k=1}^{i}\sum_{j=1}^{k} {k-1 \choose j-1}^{2} \tau^{-2(k-1)}\right) \tau^{ip(i-1)}}$$

$$= p\tau^{\frac{i(i-1)(1-p)}{2}} \sqrt{\left(\sum_{k=1}^{i}\sum_{j=1}^{k} {k-1 \choose j-1}^{2} \tau^{2(j-1)}\right) \left(\sum_{k=1}^{i}\sum_{j=1}^{k} {k-1 \choose j-1}^{2} \tau^{-2(k-1)}\right)}. \tag{3.27}$$

Remark 3.3. As is established in Theorem 3.1, the conditioning number of T increases (and so the conditioning number of Λ_i) in an exponential way with both the estimated order of the system, i, and the number of outputs, p.

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