# Nonstationary consistency of covariance based subspace methods 

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#### Abstract

In this paper we study nonstationary consistency of subspace methods for eigenstructure identification, i.e., the ability of subspace algorithms to converge to the true eigenstructure despite nonstationarities.


## I. Introduction

In 1985, Benveniste and Fuchs [6] proved that the Instrumental Variable method and what was called the Balanced Realization method for linear system eigenstructure identification are consistent for certain classes of nonstationary systems. Since this paper, the family of subspace algorithms has been invented [14], [19], [22], [23], [24] and has expanded rapidly. Therefore, we felt it was timely revisiting the results of [6] and generalizing them to subspace methods. To this end, paper [6] had first to be restructured to show up an important intermediate result, which had not been noticed explicitly. Having done this, we show how this key result can be applied in various situations.

There are a number of convergence studies on subspace methods in a stationary context in the literature, see [11], [3], [4], [5], [9], [10] to mention just a few of them. These papers provide deep and technically difficult results including convergence rates. They typically address the problem of identifying all system matrices or the transfer matrix, i.e., both the pole and zero parts of the system. In contrast, the nonstationary consistency property that we study here holds for the eigenstructure (the pole part) only and does not apply to the zero part. It is definitely different from the problem considered in [21].

The paper is organized as follows. The problem of nonstationary consistency is stated in Section II. Section III collects the key results of our approach. In Section III-A we show that nonstationary consistency follows if we can construct matrix estimators having the almost factorized form $R_{i}(N)=C A^{i} G(N)+o(1)$, where $A$ and $C$ are the state transition matrix and observation matrix for identification, and $o(1) \rightarrow 0$ when the sample size $N$ goes to infinity. Historically, subspace algorithms originated from Instrumental Variable methods; we show in Section III-C that it is in turn useful to associate instruments to subspace algorithms. Since limit theorems exist for martingales regardless of stationarity, our approach to dealing with nonstationarity is based on martingale arguments; this is treated in Section III-D. These results allow us to reformulate the design of subspace algorithms ensuring nonstationary consistency for eigenstructure estimation as the quest for suitable instruments, for linear systems; this is explained in Section III-E.

[^0]Finally, in Section IV, by using the so developed toolbox of theorems and lemmas, we prove nonstationary consistency of some representative subspace algorithms.

Due to space limitations, we focus in this paper on so-called covariance based subspace algorithms. Projection based ones are not covered. The interested reader is referred to the full paper [1] for the treatment of both families of algorithms.

## II. Problem setting

Consider the following linear system

$$
\left\{\begin{array}{l}
x_{k}=A x_{k-1}+B u_{k}+v_{k}  \tag{1}\\
y_{k}=C x_{k-1}+D u_{k}+w_{k}
\end{array}\right.
$$

where $k \in \mathbb{Z}, x$ is the $\mathbb{R}^{n}$-valued state, $u$ is the $\mathbb{R}^{m}$-valued observed input, $v$ and $w$ are unobserved input disturbances, and $y$ is the $\mathbb{R}^{q}$-valued observed output.

The key point of this work is that the unobserved input disturbances can be nonstationary. For instance, they can be white noises having unknown time-varying covariance matrices. For this case, we should rather reformulate system (1) in the following form, which enlightens that $y_{k}$ itself is nonstationary in a nontrivial way:

$$
\left\{\begin{array}{l}
x_{k}=A x_{k-1}+B u_{k}+K(k) \nu_{k}  \tag{2}\\
y_{k}=C x_{k-1}+D u_{k}+L(k) \nu_{k}
\end{array}\right.
$$

where $\left[\begin{array}{c}K(k) \\ L(k)\end{array}\right]\left[\begin{array}{ll}K^{T}(k) & L^{T}(k)\end{array}\right]$ is the time-varying covariance matrix of the excitation noise in (1), and $\nu_{k}$ is a stationary standard white noise. Note that the zero part of the transfer matrix $\left[u_{k} \nu_{k}\right] \mapsto y_{k}$ is time-varying in this case, so that consistency makes sense only w.r.t. the pole part.

The problem we consider is the identification of the pair $(C, A)$ up to a change of basis in the state space of system (2). Equivalently, we identify the pairs $\left(\lambda, C \varphi_{\lambda}\right)$, where $\lambda$ ranges over the set of eigenvalues of $A$ (the poles of system (2)) and $\varphi_{\lambda}$ are a corresponding set of eigenvectors. In words, we consider the problem of eigenstructure identification.

Our objective is to show that subspace methods provide consistent estimators of the eigenstructure, also for nonstationary cases as above. Of course, none of the matrices $A, B, C, D, K(k)$, and $L(k)$, are known. Matrices $B, D, K(k)$, and $L(k)$, are regarded as nuisance and are not for identification in this paper.

## III. BASIC THEOREMS FOR NONSTATIONARY CONSISTENCY

Throughout this paper, for $t_{N}$ a nondecreasing sequence of positive real nombers, $o\left(t_{N}\right)$ generically denotes a
matrix-valued sequence $M_{N}$, of fixed dimensions, such that $M_{N} / t_{N} \rightarrow 0$ when $N$ tends to infinity.

Also, throughout this paper, we distinguish Conditions from Assumptions. Assumptions will refer to hypothesized properties of the system or its inputs; Assumptions may or may not hold. In contrast, Conditions can be enforced; enforcing them will be typically part of the process of designing the subspace algorithms.

## A. From Hankel matrices to eigenstructure

Consider an observable pair $(C, A)$ of matrices, where $C$ is $q \times n$ and $A$ is $n \times n$. Throughout this paper, $p$ denotes an integer large enough such that

$$
\operatorname{rank}\left(\mathcal{O}_{p}\right)=n, \quad \text { where } \quad \mathcal{O}_{p} \triangleq\left[\begin{array}{l}
C  \tag{3}\\
C A \\
\vdots \\
C A^{p-1}
\end{array}\right]
$$

For $i=1, \ldots, p$ and $N>0$, consider a family $R_{i}(N)$ of $q \times r$-matrices, satisfying the following assumption:

Condition 1: The matrices $R_{i}(N), N>0$, decompose as $R_{i}(N)=C A^{i-1} G(N)+o(1)$. Furthermore, the sequence of $n \times r$-matrices $G(N), N>0$, satisfies $\liminf _{N \rightarrow \infty} \sigma_{n}(G(N))>0$, where $\sigma_{n}(M)$ denotes the $n$-th largest singular value of matrix $M$.

Generic subspace algorithm: The following generic algorithm is considered throughout this paper. It assumes a finite family of matrices $R_{i}(N)$ as above and returns a pair $(C(N), A(N))$. The sentence
" $R_{i}(N)$ provides consistent estimators for $(C, A)$ "
that we use throughout this paper means that, when provided with the sequence $R_{i}(N)$, this generic algorithm yields consistent estimators $(C(N), A(N))$ for the pair $(C, A)$, see Theorem 1 below. Consider the matrix $\mathcal{H}_{p}(N)$ defined by

$$
\mathcal{H}_{p}(N) \triangleq\left[\begin{array}{c}
R_{1}(N)  \tag{4}\\
R_{2}(N) \\
\vdots \\
R_{p}(N)
\end{array}\right]
$$

and SVD-decompose it as:

$$
\begin{aligned}
\mathcal{H}_{p}(N) & =\sum_{i=1}^{\min (p q, r)} \sigma_{i} \mathbf{u}_{i} \mathbf{v}_{i}^{T} \\
& =\sum_{i=1}^{n} \sigma_{i} \mathbf{u}_{i} \mathbf{v}_{i}^{T}+\sum_{i=n+1}^{\min (p q, r)} \sigma_{i} \mathbf{u}_{i} \mathbf{v}_{i}^{T} \\
& =\mathbf{U} \operatorname{diag}\left(\sigma_{1}, \ldots, \sigma_{n}\right) \mathbf{V}^{T}+\sum_{i=n+1}^{\min (p q, r)} \sigma_{i} \mathbf{u}_{i} \mathbf{v}_{i}^{T}
\end{aligned}
$$

Partition the above defined $p q \times n$ matrix $\mathbf{U}$ into its $p$ successive $q$-block rows $\mathbf{U}_{1}, \ldots, \mathbf{U}_{p}$ and set

$$
\mathbf{U}^{\uparrow} \triangleq\left[\begin{array}{c}
\mathbf{U}_{2} \\
\vdots \\
\mathbf{U}_{p}
\end{array}\right] \quad \text { and } \quad \mathbf{U}^{\downarrow} \triangleq\left[\begin{array}{c}
\mathbf{U}_{1} \\
\vdots \\
\mathbf{U}_{p-1}
\end{array}\right]
$$

Using these notations, set

$$
\begin{aligned}
& C(N) \triangleq \mathbf{U}_{1} \\
& A(N) \triangleq \quad \text { least-squares solution of } \mathbf{U}^{\uparrow}=\mathbf{U}^{\downarrow} A
\end{aligned}
$$

Theorem 1: [6] Under Condition 1, $(C(N), A(N))$ defined by (4-5) is a consistent estimator of $(C, A)$ in the following sense: there exists a sequence of matrices $T(N)$, with $T(N)$ and $T^{-1}(N)$ uniformly bounded w.r.t. $N$, such that $\lim _{N \rightarrow \infty} T^{-1}(N) A(N) T(N) \rightarrow A$, and $\lim _{N \rightarrow \infty} C(N) T(N) \rightarrow C$.

Proof: It is found in [6], second part of Section III-C, dealing with the Balanced Realization algorithm. Besides the fact that reference [6] speaks $(H, F, G)$ and not $(A, B, C)$, the only slight change is that matrix $G(N)$ in Condition (1) replaces the controllability matrix $\mathbf{C}\left(F, G_{S}\right)$ of [6], where $S$ is the sample length.

## B. Notations

For $X$ and $Y$ two matrices of compatible dimensions, define:

$$
\begin{align*}
\langle X, Y\rangle & \triangleq X Y^{T} \\
\|X\|^{2} & \triangleq \operatorname{Tr}\langle X, X\rangle \\
\mathrm{E}(X \mid Y) & \triangleq\langle X, Y\rangle\langle Y, Y\rangle^{\dagger} Y  \tag{6}\\
\mathrm{E}\left(X \mid Y^{\perp}\right) & \triangleq X-\mathrm{E}(X \mid Y)
\end{align*}
$$

where Tr denotes the trace and superscript ${ }^{\dagger}$ denotes the pseudo inverse. For $\left(y_{k}\right)_{k \in \mathbb{Z}}$ a $\mathbb{R}^{q}$-valued data sequence and $N>0$ a window length, define

$$
Y_{i}(N) \triangleq \triangleq\left[\begin{array}{llll}
y_{i+N-1} & \ldots & y_{i+1} & y_{i}
\end{array}\right]
$$

and write simply $Y_{i}$ if $N$ is understood. For $\left(x_{k}\right)_{k \in \mathbb{Z}}$ and $\left(z_{k}\right)_{k \in \mathbb{Z}}$ two data sequences of compatible dimensions, we write:

$$
\left\langle X_{i}, Z_{j}\right\rangle_{N} \triangleq\left\langle X_{i}(N), Z_{j}(N)\right\rangle
$$

and

$$
\mathrm{E}_{N}\left(X_{i} \mid Z_{j}\right) \triangleq \mathrm{E}\left(X_{i}(N) \mid Z_{j}(N)\right)
$$

Finally, we shall make use of the following data Hankel matrices:

$$
\mathcal{Y}_{i, M}^{+}(N) \triangleq\left[\begin{array}{c}
Y_{i+M} \\
\vdots \\
Y_{i+2} \\
Y_{i+1}
\end{array}\right], \mathcal{Y}_{i, M}^{-}(N) \triangleq\left[\begin{array}{c}
Y_{i} \\
Y_{i-1} \\
\vdots \\
Y_{i-M}
\end{array}\right]
$$

and

$$
\mathcal{Y}_{i, M}(N) \triangleq\left[\begin{array}{l}
\mathcal{Y}_{i, M}^{+} \\
\mathcal{Y}_{i, M}^{-}
\end{array}\right]
$$

The above notations are introduced because, depending on the considered algorithms, the data set is indexed as $y_{N}, \ldots, y_{1}$ (when only "future" data are needed), or $y_{N}, \ldots, y_{1}, y_{0}, \ldots, y_{-N}$ (when data are split into future and past). Many authors use rather $y_{1}, \ldots, y_{N}, y_{N+1}, \ldots, y_{2 N}$, or variants thereof. Clearly, the difference is only notational. Also, we have taken identical index $M$ in $\mathcal{Y}_{i, M}^{+}$and $\mathcal{Y}_{i, M}^{-}$ when building $\mathcal{Y}_{i, M}$. Of course, we could take different indices $M_{+}$and $M_{-}$whithout impairing the validity of what follows.

Finally in order to refer to the different algorithms in a systematic way in the sequel, we shall superscript the referred $R_{i}(N)$ with the index of the corresponding equation. For example,

$$
\begin{equation*}
R_{i}^{(12)}(N) \text { denotes } R_{i}(N) \text { as specified by (12). } \tag{7}
\end{equation*}
$$

## C. Instruments

In this section, we revisit the old concept of "instrument" and use it in our context. Unlike in Section II where our problem was stated, we do not distinguish here between observed and unobserved inputs. In the following system, vector $\xi$ collects all inputs of the system considered throughout this section:

$$
\left\{\begin{array}{l}
x_{k}=A x_{k-1}+B^{\prime} \xi_{k}  \tag{8}\\
y_{k}=C x_{k-1}+D^{\prime} \xi_{k}
\end{array}\right.
$$

In (8), $k \in \mathbb{Z}, x$ is the $\mathbb{R}^{n}$-valued state, $\xi$ is the $\mathbb{R}^{m}$-valued input, and $y$ is the $\mathbb{R}^{q}$-valued observed output. Fix a window length $N$. With the notations of Section III-B, system (8) rewrites as follows, for $i=1, \ldots, p$ :

$$
\left\{\begin{align*}
X_{i} & =A X_{i-1}+B^{\prime} \Xi_{i}  \tag{9}\\
Y_{i} & =C X_{i-1}+D^{\prime} \Xi_{i}
\end{align*}\right.
$$

In the following lemma we introduce instruments as the key tool in our analysis:

Lemma 1: (instruments) Let $\left(z_{k}\right)_{k \in \mathbb{Z}}$ be an $\mathbb{R}^{M}$-valued data sequence and $\left(s_{N}\right)_{N>0}$ an $\mathbb{R}_{+}$-valued sequence such that

$$
\text { for } \begin{align*}
j \in & \{1, \ldots, i\}:\left\langle\Xi_{j}, Z_{0}\right\rangle_{N}=o\left(s_{N}\right)  \tag{10}\\
& \liminf _{N \rightarrow \infty} \sigma_{n}\left(\frac{1}{s_{N}}\left\langle X_{0}, Z_{0}\right\rangle_{N}\right)>0 \tag{11}
\end{align*}
$$

Then,

$$
\begin{equation*}
R_{i}(N) \triangleq \frac{1}{s_{N}}\left\langle Y_{i}, Z_{0}\right\rangle_{N} \tag{12}
\end{equation*}
$$

satisfies Condition 1. In the sequel, we call instrument $a$ signal $\left(z_{k}\right)$ satisfying (10) and (11) for some sequence $s_{N}$.

Proof: The following decompositions hold, for $i>0$ :

$$
y_{k+i}=C A^{i-1} x_{k}+\sum_{j=1}^{i-1} C A^{i-1-j} B^{\prime} \xi_{k+j}+D^{\prime} \xi_{k+i}
$$

with the convention that $\sum_{1}^{0}=0$, and

$$
\begin{align*}
& \sum_{k=0}^{N-1} y_{k+i} z_{k}^{T}= \\
& C A^{i-1} \sum_{k=0}^{N-1} x_{k} z_{k}^{T}+  \tag{13}\\
& \sum_{j=1}^{i-1} C A^{i-1-j} \sum_{k=0}^{N-1} B^{\prime} \xi_{k+j} z_{k}^{T}+\sum_{k=0}^{N-1} D^{\prime} \xi_{k+i} z_{k}^{T}
\end{align*}
$$

Equation (13) rewrites as follows:

$$
\begin{align*}
& \left\langle Y_{i}, Z_{0}\right\rangle_{N} \\
& =C A^{i-1}\left\langle X_{0}, Z_{0}\right\rangle_{N}+  \tag{14}\\
& \sum_{j=1}^{i-1} C A^{i-1-j} B^{\prime}\left\langle\Xi_{j}, Z_{0}\right\rangle_{N}+D^{\prime}\left\langle\Xi_{i}, Z_{0}\right\rangle_{N}
\end{align*}
$$

which proves that $R^{(12)}(N)$ satisfies Condition 1, thanks to (10) and (11).

Lemma 1 and Theorem 1 together ensure that $R^{(12)}(N)$ provides consistent estimators for the pair $(C, A)$ (see (7) for the notational convention used here).

Applying Lemma 1 to system (1) with its combined observed and unobserved inputs can be (tentatively) performed via the following substitutions:

$$
\left[\begin{array}{c}
B^{\prime}  \tag{15}\\
D^{\prime}
\end{array}\right] \xi_{k}=\left[\begin{array}{l}
B \\
D
\end{array}\right] u_{k}+\left[\begin{array}{c}
v_{k} \\
w_{k}
\end{array}\right]
$$

Of course, if input $\xi_{k}$ is observed, i.e., $v_{k}=w_{k}=0$ in (15), then one can chose instrument $z_{k}$ in such a way that $\left\langle\Xi_{j}, Z_{0}\right\rangle_{N}=0$ exactly. This is no longer feasible if unobserved inputs exist, since $\Xi_{j}$ is no longer observed in this case. Therefore, additional work is needed for analysing system (1) with its combined observed/unobserved inputs. The next section addresses this missing point.

## D. A martingale argument

Let us discuss the key conditions allowing us to apply Lemma 1 and Theorem 1 to system (1), taking the unobserved inputs $v$ and $w$ into account.

Suppose first that there is no unobserved input disturbance, i.e., $v=w=0$ in (1). Then, the observed values for the $U_{i}$ 's can be used while constructing instrument $z_{k}$ in Lemma 1 , since these are known-this is exactly what happens in the so-called "deterministic" subspace methods [19], where $\Xi_{j}=U_{j}$ holds and $\left\langle U_{j}, Z_{0}\right\rangle_{N}=0$ is enforced. Note that no assumption of stochastic nature is required for this reasoning.

Next, consider the opposite case in which there is no observed input, i.e., $B=D=0$ in (1). Since input disturbances are not observed, the actual values of $\Xi_{j}$ are unknown when applying Lemma 1 and therefore cannot be used while constructing the instrument $z_{k}$.

This problem, however, can be solved by using stochastic knowledge about unobserved input disturbances. To this end, we now introduce the needed probabilistic setting.

Lemma 2: Let $\left(v_{k}\right)_{k \geq 0}$ and $\left(z_{k}\right)_{k \geq 0}$ be two sequences of square integrable vector valued random variables defined over some probability space $(\Omega, \mathcal{G}, \mathbb{P})$ and let $\left(\mathcal{G}_{k}\right)_{k \geq 0}$ be an increasing family of sub- $\sigma$-algebras of $\mathcal{G}$ such that:

$$
\begin{array}{r}
\sup _{k \geq 0} \mathbb{E}\left\|v_{k}\right\|^{2} \leq K<\infty \\
\lim _{N \rightarrow \infty} \sum_{k=0}^{N \geq}\left\|z_{k}\right\|^{2}=+\infty \text { w.p. } 1  \tag{16}\\
v_{k} \text { and } z_{k} \text { are } \mathcal{G}_{k}-\text { measurable, and } \\
\mathbb{E}\left(v_{k} \mid \mathcal{G}_{k-1}\right)=0
\end{array}
$$

Then $\lim _{N \rightarrow \infty} \frac{M_{N}}{\sum_{k=0}^{N}\left\|z_{k}\right\|^{2}}=0$ holds w.p. 1, where $M_{N} \triangleq$ $\sum_{k=j}^{N} v_{k} z_{k-j}^{T}$ for $j>0$.

Nota: In formula (16), the conditional expectation $\mathbb{E}\left(. \mid \mathcal{G}_{k-1}\right)$ should not be confused with our matrix projection operator $\mathrm{E}(. \mid$.$) in (6).$

Proof: It is a mild variation of the argument of [6], Section III-A. We repeat it here for the sake of completeness. Since we can reason on each entry of matrix $M_{N}$ separately, we can, without loss of generality, assume that $v_{k}$ and $z_{k}$ are both scalar signals. By the second condition of (16), we know that $\left(M_{k}\right)_{k \geq 0}$ is a square integrable scalar martingale w.r.t. $\left(\mathcal{G}_{k}\right)_{k \geq 0}$. By (16), we have $\mathbb{E}\left(\left(M_{k}-M_{k-1}\right)^{2} \mid \mathcal{G}_{k-1}\right)=$ $\mathbb{E}\left(v_{k}^{2} \backslash \mathcal{G}_{k-1}\right) z_{k-j}^{2}=\mathbb{E}\left(v_{k}^{2}\right) z_{k-j}^{2} \leq K z_{k-j}^{2}$. The proof is
then completed by using Theorem 2 below, which can be found in [13], [17].

Theorem 2: ([13], [17]) Let $\left(M_{k}\right)_{k \geq 0}$ be a locally square integrable martingale w.r.t. $\left(\mathcal{G}_{k}\right)_{k \geq 0}$, such that $M_{0}=0$. Set

$$
[M, M]_{k}=\sum_{l=1}^{k} \mathbb{E}\left(\left(M_{l}-M_{l-1}\right)^{2} \mid \mathcal{G}_{l-1}\right)
$$

Then, w.p.1, we have $\frac{M_{k}}{[M, M]_{k}} \rightarrow 0$ on the set $\left\{\lim _{k \rightarrow \infty}[M, M]_{k}=+\infty\right\}$, and $\lim _{k \rightarrow \infty} M_{k}$ exists and is finite on the set $\left\{\lim _{k \rightarrow \infty}[M, M]_{k}<+\infty\right\}$.

## E. Analysing covariance based algorithms

In this section we combine the results from Sections III-C and III-D to handle system (1) with its combined observed/unobserved inputs. We repeat again system (1) for convenience:

$$
\left\{\begin{array}{l}
x_{k}=A x_{k-1}+B u_{k}+v_{k}  \tag{17}\\
y_{k}=C x_{k-1}+D u_{k}+w_{k}
\end{array}\right.
$$

where $k \in \mathbb{Z}, y$ is the $\mathbb{R}^{q}$-valued observed output, $x$ is the $\mathbb{R}^{n}$-valued state, $u$ is the $\mathbb{R}^{m}$-valued observed input, $(v, w)$ is an unobserved input disturbance.

To be able to use stochastic information on the unobserved inputs $v, w$ we assume that all variables arising in system (17) are defined over some probability space $(\Omega, \mathcal{F}, \mathbb{P})$.

Available information is captured by the following $\sigma$ algebras:

$$
\begin{aligned}
& \mathcal{F}_{k} \triangleq \underbrace{\sigma\left(u_{j}: j \in \mathbb{Z}\right)}_{\mathcal{F} u} \vee \underbrace{\sigma\left(y_{l}, v_{l}, w_{l}: l \leq k\right)}_{\mathcal{F}_{k}^{\prime, v, w}} \\
& \mathcal{F}_{k}^{o} \triangleq \underbrace{\sigma\left(u_{j}: j \in \mathbb{Z}\right)}_{\mathcal{F}^{u}} \vee \underbrace{\sigma\left(y_{l}: l \leq k\right)}_{\mathcal{F}_{k}^{y}}
\end{aligned}
$$

$\sigma$-algebra $\mathcal{F}^{u}$ is the information provided by the entire observed input sample; $\sigma$-algebra $\mathcal{F}_{k}^{y, v, w}$ is the information provided by the unobserved inputs $v$ and $w$ and the output $y$ up to time $k$; finally, $\sigma$-algebra $\mathcal{F}_{k}^{y}$ is the information provided by the only output $y$ up to time $k$. Regarding the unobserved inputs, we assume the following:

Assumption 1: [unobserved inputs] Stochastic inputs $v$ and $w$ satisfy the following conditions:

$$
\begin{aligned}
& \sup _{k \geq 0} \mathbb{E}\left(\left\|v_{k}\right\|^{2}+\left\|w_{k}\right\|^{2}\right)<\infty, \text { and for } j>0: \\
& \forall k \geq 0, \mathbb{E}\left(v_{k+j} \mid \mathcal{F}_{k}\right)=0 \text { and } \mathbb{E}\left(w_{k+j} \mid \mathcal{F}_{k}\right)=0
\end{aligned}
$$

Note that these conditions do not request any kind of stationarity. Assumption 1 involves the joint distribution of $v_{k}, w_{k}$, and $u_{k}$. It is in particular satisfied when observed and unobserved inputs are independent. Besides Assumption 1, no condition is required on the statistics of the observed input $u_{k}$. Consider the following conditions:

Condition 2: (instruments) Instrument $\left(z_{k}\right)$ satisfies the following conditions:

$$
\begin{array}{r}
z_{k} \text { is } \mathcal{F}_{k}^{o} \text {-measurable } \\
\lim _{N \rightarrow \infty} s_{N}=\infty, \text { where } s_{N} \triangleq \sum_{k=-M}^{N-1}\left\|z_{k}\right\|^{2} \\
\left\langle\left[\begin{array}{c}
B \\
D
\end{array}\right] U_{j}, Z_{0}\right\rangle_{N}=o\left(s_{N}\right) \text { for } j>0 \\
\liminf _{N \rightarrow \infty} \sigma_{n}\left(\frac{1}{s_{N}}\left\langle X_{0}, Z_{0}\right\rangle_{N}\right)>0 \tag{21}
\end{array}
$$

Property (18) guarantees that instrument $z_{k}$ depends only on observed quantities. Integer $M \geq 0$ in (19) is a constant selected according to each particular instance of the family $R_{i}(N)$. Property (19) expresses that instrument $\left(z_{k}\right)$ possesses sustained energy. The following theorem is our first main result. It provides the analysis of algorithms of the form (12), i.e., covariance based ones.

Theorem 3: (covariance based subspace) Assume that Assumption 1 regarding unobserved inputs, and Condition 2 regarding instruments, are in force. Then, $R_{i}^{(12)}(N)$ satisfies Condition 1.

Therefore, $R_{i}^{(12)}(N)$ yields a consistent subspace algorithm, by Theorem 1 .

Proof: Using the notations of Section III-C, system (17) writes as follows, for $i=1, \ldots, p$ :

$$
\left\{\begin{align*}
X_{i} & =A X_{i-1}+B U_{i}+V_{i}  \tag{22}\\
Y_{i} & =C X_{i-1}+D U_{i}+W_{i}
\end{align*}\right.
$$

On the other hand, system (17) yields the following decomposition for $y_{k+i}, i>0$ (we use the convention that $\sum_{1}^{0}=0$ ):

$$
y_{k+i}=C A^{i-1} x_{k}+\sum_{j=1}^{i-1} C A^{i-1-j} \widehat{v}_{k+j}+\widehat{w}_{k+i}
$$

where $\widehat{v}_{k} \triangleq B u_{k}+v_{k}$ and $\widehat{w}_{k} \triangleq D u_{k}+w_{k}$. Using the notations of Section III-C, this decomposition rewrites as follows, for $i>0$ :

$$
\begin{equation*}
Y_{i}=C A^{i-1} X_{0}+\sum_{j=1}^{i-1} C A^{i-1-j} \widehat{V}_{j}+\widehat{W}_{i} \tag{23}
\end{equation*}
$$

where $\widehat{V}_{i} \triangleq B U_{i}+V_{i}$ and $\widehat{W}_{i} \triangleq D U_{i}+W_{i}$. Note that

$$
\begin{equation*}
\left\langle V_{j}, Z_{0}\right\rangle_{N}=\sum_{k=0}^{N-1} v_{k+j} z_{k}^{T} \tag{24}
\end{equation*}
$$

and a similar formula holds with $W_{j}$ instead of $V_{j}$. By (18) and (19) of Condition 2, instrument $\left(z_{k}\right)$ satisfies (16) in Lemma 2. By Assumption 1, noises $v_{k}$ and $w_{k}$ satisfy (16) in Lemma 2, with $\mathcal{F}_{k}$ substituted for $\mathcal{G}_{k}$. Therefore Lemma 2 can be applied with $\mathcal{F}_{k}$ substituted for $\mathcal{G}_{k}$, which yields, for any $j \in\{1, \ldots, p\}$ :

$$
\left\langle\left[\begin{array}{l}
V_{j}  \tag{25}\\
W_{j}
\end{array}\right], Z_{0}\right\rangle_{N}=o\left(s_{N}\right)
$$

Set

$$
\begin{aligned}
\xi_{k} & \triangleq\left[\begin{array}{c}
B \\
D
\end{array}\right] u_{k}+\left[\begin{array}{c}
v_{k} \\
w_{k}
\end{array}\right] \\
B^{\prime} & \triangleq\left[\begin{array}{ll}
I_{n} & 0_{q}
\end{array}\right] \\
D^{\prime} & \triangleq\left[\begin{array}{ll}
0_{n} & I_{q}
\end{array}\right]
\end{aligned}
$$

where the subscripts ${ }_{n}$ and ${ }_{q}$ indicate the dimensions of the corresponding matrices. Using this change of notation allows us to rewrite system (17) in the form (8) used in Lemma 1.

Consider now Condition 2. Combining (25) with (20) shows that system (8) satisfies (10) in Lemma 1. On the other hand (11) in Lemma 1 is ensured by Property (21) of Condition 2. Therefore, by Lemma 1 we conclude that Condition 1 is satisfied.

Remarks: 1/ In fact our method could accomodate as well additional "small" perturbations in system (17), i.e., additional inputs $\mu_{k}$ and $\nu_{k}$ in state and observation equations respectively, such that $\frac{1}{s_{N}} \sum_{k=-M}^{N-1}\left\|\mu_{k}\right\|^{2}+\left\|\nu_{k}\right\|^{2}=o(1)$. Transient terms or leakage effects such as considered in [7], [8] are covered by these additional terms, and therefore do not impair nonstationary consistency.

2/ Projection based subspace methods, i.e., methods of the form

$$
\begin{equation*}
R_{i}(N) \triangleq \frac{1}{s_{N}} \mathrm{E}_{N}\left(Y_{i} \mid Z_{0}\right) \tag{26}
\end{equation*}
$$

are in fact more popular than covariance based ones, see [20]. They are often referred to as "data based" subspace methods. Unfortunately, these methods cannot be handled directly by Theorem 3. In fact, Theorem 1 itself does not apply. The reason for this is simple: $R_{i}^{(26)}(N)$ has dimensions $q \times$ $\operatorname{dim}\left(Z_{0}(N)\right)$. So its dimensions vary with $N$ and therefore Theorem 1 cannot apply. In the full paper [1] a weighting technique is used to overcome this difficulty.

## IV. AnAlysis of some subspace algorithms

In this section we apply our toolbox of theorems and lemmas to sample subspace methods. To avoid boring notational adjustments, we keep our notational conventions and will therefore sometimes deviate from the original presentations in this respect.

Key conditions ensuring nonstationary consistency are Assumption 1 and Condition 2. Assumption 1 involves the unobserved inputs, we assume it to hold throughout this section and will not discuss it any further. In contrast, Condition 2 is a design constraint on the selection of the instruments: this is the key condition to be verified or enforced when analysing specific algorithms.

Regarding the details of Condition 2, we shall pay great attention to verifying that (18) and (20) are satisfied, as these conditions drive the choice of the instruments. Condition (21) amounts to requiring that the instrument is well correlated to the state. In contrast, we shall not discuss the satisfaction of condition (19); this condition just translates, for each particular algorithm, into corresponding conditions for the original system (17).

Finally, checking for consistency requires that proper normalization is applied. This is the very role of the scaling
factor $s_{N}$. In practice the algorithms are applied with given sample length $N$, and then, scaling is just an irrelevant issue. Therefore, we shall ignore scaling in this section.

## A. Output-only (OO) subspace algorithms

By definition, these algorithms assume $B=D=0$ in (17). Therefore (20) in Condition 2 is trivially satisfied, thus we essentially need to check the measurability property (18).

Basic OO subspace algorithm: This is the simplest algorithm to analyse. Introduce the instrument

$$
z_{k} \triangleq\left[\begin{array}{c}
y_{k}  \tag{27}\\
\vdots \\
y_{k-M}
\end{array}\right]
$$

and take

$$
\begin{equation*}
R_{i}(N)=\left\langle Y_{i}, Z_{0}\right\rangle_{N} \tag{28}
\end{equation*}
$$

Instrument (27) satisfies (18) in Condition 2. Hence Theorem 3 applies and proves consistency of $R_{i}^{(28)}(N)$. Note that

$$
\left\langle X_{0}, Z_{0}\right\rangle_{N}=\left[\begin{array}{llll}
F(N) & A F(N) & \ldots & A^{M} F(N)
\end{array}\right]
$$

where $F(N)=\left\langle X_{0}, Y_{0}\right\rangle_{N}$. Hence, (21) can be interpreted as $y_{k}$ being "uniformly of order $n$ ".

Covariance driven OO subspace algorithm [6], [18], [2], [12]: This algorithm is a variation of the previous algorithm, it was however proposed earlier. It consists in computing, for $i=1, \ldots, p$ :

$$
\begin{align*}
& R_{i}(N)=\left[\begin{array}{llll}
\hat{r}_{i}(N) & \hat{r}_{i+1}(N) & \cdots & \hat{r}_{i+M}(N)
\end{array}\right]  \tag{29}\\
& \text { where } \hat{r}_{j}(N)=\left\langle Y_{j}, Y_{0}\right\rangle_{N}
\end{align*}
$$

With instrument $z_{k}$ as in (27), we have

$$
R_{i}^{(29)}(N)-\left\langle Y_{i}, Z_{0}\right\rangle_{N}=\left[\begin{array}{lll}
\delta\left\langle Y_{i}, Y_{0}\right\rangle & \ldots & \delta\left\langle Y_{i}, Y_{-M}\right\rangle
\end{array}\right]
$$

where $\delta\left\langle Y_{i}, Y_{-k}\right\rangle \triangleq\left\langle Y_{i}, Y_{-k}\right\rangle_{N}-\left\langle Y_{i+k}, Y_{0}\right\rangle_{N}$ is such that

$$
\left\|\delta\left\langle Y_{i}, Y_{-k}\right\rangle\right\| \leq 2 s_{M, N}^{*}
$$

where

$$
s_{M, N}^{*} \triangleq \sup _{-M \leq j \leq N-M} \sum_{l=j}^{j+M}\left\|y_{l}\right\|^{2}
$$

This implies

$$
\begin{equation*}
\left\|R_{i}^{(29)}(N)-\left\langle Y_{i}, Z_{0}\right\rangle_{N}\right\|=o\left(s_{N}\right) \tag{30}
\end{equation*}
$$

provided that the following assumption holds:
Assumption 2: For $M$ fixed, $s_{M, N}^{*}=o\left(s_{N}\right)$.
With Assumption 2, (30) holds and therefore instrument $z_{k}$ defined in (27) satisfies Condition 2. By Theorem 3, we derive that $R_{i}^{(29)}(N)$ yields a consistent subspace algorithm.

## B. Input-output (IO) subspace algorithms

We review some representative variants.

Covariance driven IO subspace algorithms with projection on the orthogonal of the input [12]: This algorithm consists in computing (cf. notations (6)):

$$
\begin{aligned}
\hat{r}_{i}(N) & =\left\langle Z_{i}, Z_{0}\right\rangle_{N}, \text { where } Z_{i} \triangleq \mathrm{E}_{N}\left(Y_{i} \mid \mathcal{U}_{0, M}^{\perp}\right) \\
R_{i}(N) & =\left[\begin{array}{llll}
\hat{r}_{i}(N) & \hat{r}_{i+1}(N) & \ldots & \hat{r}_{i+M}(N)
\end{array}\right](31)
\end{aligned}
$$

First, note that $\left\langle Z_{i}, Z_{0}\right\rangle_{N}=\left\langle Y_{i}, Z_{0}\right\rangle_{N}$. The associated instrument $z_{k}$ is therefore the sequence of the successive columns of matrix $Z_{0}$. Note that $z_{k}$ is $\mathcal{F}_{k}^{o}$-measurable. The rest of the analysis of this algorithm proceeds as for $R_{i}^{(29)}(N)$. Property (21) can be seen as that $z_{k}$ itself being "uniformly of order $n$ ".

Covariance driven subspace algorithm using projected past inputs and outputs as instruments [25]: Those methods encompass the methods also known as IVM, CVA, POMOESP and N4SID in their covariance form [25]. In this paper, we will focus on the unweighted IV related to $\mathcal{H}_{p}$ defined as

$$
\begin{equation*}
\mathcal{H}_{p}=\left\langle\mathcal{Y}_{0, p}^{+}, \mathcal{L}_{0, M}^{-}\right\rangle_{N} \tag{32}
\end{equation*}
$$

where $\mathcal{L}_{0, M}^{-}$is defined by stacking for $i=-M, \ldots, 0$

$$
L_{i} \triangleq \mathrm{E}_{N}\left(\left.\left[\begin{array}{c}
U_{i}  \tag{33}\\
Y_{i}
\end{array}\right] \right\rvert\,\left(\mathcal{U}_{0, M}^{+}\right)^{\perp}\right)
$$

The rest of the analysis of this algorithm proceeds as for (31).

## C. Time-vs. frequency-domain.

For $\left(y_{k}\right)_{k \in \mathbb{Z}}$ an $\mathbb{R}^{q}$-valued data sequence and $N_{\widehat{\prime}}>0$ a window length, the DFT of $Y_{i}(N)$, denoted by $\widehat{Y}_{i}(N)$, is equal to $\widehat{Y}_{i}(N)=Y_{i}(N) \Delta_{N}^{q}$, where $\Delta_{N}^{q}$ is the $N$ DFT matrix. Since matrix $\Delta_{N}^{q}$ is orthogonal, $\langle\widehat{X}, \widehat{Y}\rangle_{N}=$ $\langle X, Y\rangle_{N}$, and $\mathrm{E}_{N}(\widehat{X} \mid \widehat{Y})=\mathrm{E}_{N}(X \mid Y) \Delta_{N}^{q}$. Hence, Condition 2 can be considered equivalently in the time domain or in the frequency domain. Therefore, frequency domain subspace algorithms corresponding to [15], [16] behave exactly the same way as their time domain counterparts regarding nonstationary consistency.

## V. Conclusion

We have revisited eigenstructure identification via subspace methods. We have shown that consistency still holds in case of nonstationary inputs (in fact, for "nonstationary zero part").

Martingale techniques were used to deal with unobserved inputs-for these, "deterministic" projections based on observed data cannot be used; they can however be supplemented by "stochastic" projections using conditional expectations.

Not surprisingly, transient and leakage effects are not an issue for nonstationary consistency. And the results equivalently apply to both time- and frequency-domain methods.

The simple techniques developed in this paper are not able of addressing projection based algorithms. The interested reader is referred to the full version [1], where all classes of subspace methods are encompassed.

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