# Large Sample Approximations for Variance-Covariance Matrices of High-Dimensional Time Series 

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## Outline

(1) Introduction and Overview of Actitivities Related to Data Science
(2) Projection-Based Analysis
(3) Framework and Assumptions
(9) Large Sample Approximation
(6) Applications

## Introduction

## Big Data: Sources

- Internet, sensors, cameras, simulations, ...


## Aims:

- Extract information, 'knowledge'
- Build predictive models
- Simulate scenarios
- Separate 'structure' and 'noise'
- ...

Complex and high-dimensional data ('big data')

- Functional Data (discretely observed processes, measurement curves, ...)
- Image data, video data
- High-dimensional (vector) correlated data
- Time series structure (temporal correlations)


## Overview: Stochastics Group

## Overview of Activities Related to Data Science

Questions we address:

- How to monitor (possibly high-dimensional) data streams?
- How to monitor image streams?
- How to analyze spatial-temporal correlated data?
- How to analyze high-dimensional highly-correlated vector time series? $\rightarrow$ focus of this talk.


## Overview: Stochastics Group

## Image Data (I)

Example: Preprocessing \& analysis of electroluminescence images of solar panels.
PVStatLab-Project: PV-Scan (with TÜV Rheinland, ISC Konstanz, Wrocław UoT, BMWi funded), http://www.pvstatlab.rwth-aachen.de


Figure: Example: Preprocessing using robust regression

## Overview: Stochastics Group

## Image Data (II)

Example: Image Analysis. Image as a random field $\left\{\xi_{i j}\right\}$.


$$
\begin{aligned}
& \qquad \begin{array}{c}
H_{0}: \mathrm{E}\left(\xi_{i j}\right)=\mathrm{E}\left(\xi_{u v}\right), \\
\text { for all }(i, j) \in C,(u, v) \in D \\
\\
H_{1}: \mathrm{E}\left(\xi_{i j}\right) \neq \mathrm{E}\left(\xi_{u v}\right), \\
\text { for all }(i, j) \in C,(u, v) \in D
\end{array}
\end{aligned}
$$

Figure: Regions $C$ and $D$
Aim: Asymptotic significance test taking into account spatial correlations (ongoing work), detection of defects.

## Overview: Stochastics Group

## Recent Related Publications:

1. Sovetkin, E. and Steland, A. (2015). On statistical preprocessing of PV field image data using robust regression. In:
N. E. Mastorakis, A. Ding \& M. V. Shitikova, Advances in Mathematics and Statistical Sciences, Vol. 40.
2. Steland, A. (2015). Vertically weighted averages in Hilbert spaces and applications to imaging: Fixed sample asymptotics and efficient sequential two-stage estimation, Sequential Analysis, 34 (3), 295-323.

## Overview: Stochastics Group

## Monitoring of Multivariate Data and Image Streams

- Aim: Nonparametric detection of changes
- Observe discretely sampled function representing the true signal(s) resp. image(s)
- Approaches: Hilbert-space valued r.e., random fields, Shannon/Whittaker


## Recent Related publications:

1. Prause, A. and Steland, A. (2015). Detecting changes in spatial-temporal image data based on quadratic forms. In: Stochastic Models, Statistics and Their Applications, 139-147.
2. Prause, A. and Steland, A. (2015). Sequential detection of three-dimensional signals under dependent noise, submitted.
3. Prause, A. (2015). Ph. D. thesis (finished)

## Introduction

## Large-Sample Approximations of High-Dimensional Vector Time Series

- project with R. v. Sachs (since $11 / 2013$ )
- new DFG project just started


## Setting:

Massive data set with observations on a large number of variables (features).
Focus: Analyze Dependencies

## Introduction

High-dimensional variance-covariance matrices play a crucial role in those areas, since they provide information on the dependence of the coordinates (2nd order).

The sample covariance matrix is regarded a poor estimator, since it is not consistent w.r.t. to the operator norm if the dimension is larger than the sample size $(d / n \rightarrow c>0)$.
Previous works: Banding/tapering (Bickel \& Levina, 2008), Thresholding (Chen at al., 2013), Shrinkage (Böhm and v. Sachs, 2009), ...

## Introduction

Basic problem: Observe a large number, $d=d_{n}$, of variables, $n$ repetitions (over time).

## Preliminary data analyses (preprocessing):

Frequently, e.g. by preprocessing methods, one may classify the variables in (a small number of) groups, such that

- the within-group correlation is high but
- the between-group correlation is low/negligible.

We are faced with the problem to model and analyze high-dimensional data for highly correlated variables.

## Projection-Based Analysis

Observe $d=d_{n}$ time series

$$
Y_{i}^{(\nu)}, \ldots, Y_{i}^{(\nu)}, \quad \nu=1, \ldots, d, 1 \leq i \leq n
$$

This means, we are given a vector time series of length $n$,

$$
\mathbf{Y}_{n i}=\left(Y_{i}^{(1)}, \ldots, Y_{i}^{\left(d_{n}\right)}\right)^{\prime}, \quad 1 \leq i \leq n
$$

of dimension $d_{n}$, constituting the $\left(n \times d_{n}\right)$-dimensional data matrix

$$
\mathcal{Y}_{n}=\left(Y_{i}^{(j)}\right)_{1 \leq i \leq n, 1 \leq j \leq d_{n}}
$$

We focus on second moments and thus assume $E\left(Y_{i}^{(j)}\right)=0$ for all $i=1, \ldots, n$ and $j=1, \ldots, d_{n}$.

## Projection-Based Analysis

Assume for a moment that $\mathbf{Y}_{n 1}, \ldots, \mathbf{Y}_{n n}$ is stationary. Generic copy:

$$
\mathbf{Y}_{n}=\left(Y^{(1)}, \ldots, Y^{\left(d_{n}\right)}\right)^{\prime}
$$

Unknown $\left(d_{n} \times d_{n}\right)$-dimensional sample variance-covariance matrix

$$
\boldsymbol{\Sigma}_{n}=E\left(\mathbf{Y}_{n} \mathbf{Y}_{n}^{\prime}\right)=\left(E\left(Y^{(\nu)} Y^{(\mu)}\right)\right)_{1 \leq \nu, \mu \leq d_{n}}
$$

Sample variance-covariance matrix

$$
\begin{equation*}
\widehat{\boldsymbol{\Sigma}}_{n}=\frac{1}{n} \sum_{i=1}^{n} \mathbf{Y}_{n i} \mathbf{Y}_{n i}^{\prime}=\frac{1}{n} \mathcal{Y}_{n}^{\prime} \mathcal{Y}_{n} \tag{1}
\end{equation*}
$$

Unpleasant properties for $d_{n} \gg n$, when studied as a matrix-valued estimator of $\boldsymbol{\Sigma}_{n}$, i.e. in dimension $d_{n} \times d_{n}$

## Projection-Based Analysis

But typically, one is interested in a (set of) linear combination(s) $\mathbf{w}_{n}^{\prime} \mathbf{Y}_{n}$ of the coordinates. Consider projections

$$
T_{n}=\mathbf{w}_{n}^{\prime} \mathbf{Y}_{n}
$$

for weighting vectors

$$
\mathbf{w}_{n}=\left(w_{1}, \ldots, w_{d_{n}}\right)^{\prime}, \quad n \geq 1
$$

of weights $w_{j}=w_{d_{n} j}$, not necessarily non-negative, with

$$
\begin{equation*}
\sup _{n \in \mathbb{N}}\left\|\mathbf{w}_{n}\right\|_{\ell_{1}}=\sup _{n \in \mathbb{N}} \sum_{\nu=1}^{d_{n}}\left|w_{j}\right|<\infty \tag{2}
\end{equation*}
$$

Amongst others (later), such projections allow to study single covariances between coordinates.

## Projection-Based Analysis

The projection $\mathbf{w}_{n}^{\prime} \mathbf{Y}_{n}$ has variance $\mathbf{w}_{n}^{\prime} \boldsymbol{\Sigma}_{n} \mathbf{w}_{n}$. Canonical estimator

$$
\widehat{\operatorname{Var}}\left(\mathbf{w}_{n}^{\prime} \mathbf{Y}_{n}\right)=\mathbf{w}_{n}^{\prime} \widehat{\boldsymbol{\Sigma}}_{n} \mathbf{w}_{n}
$$

behaves well for weighting vectors which select a finite number of coordinates.

Change-point problem: Test for a change in the variance of such a projection,

$$
\sigma_{n}^{2}(i)=\operatorname{Var}\left(\mathbf{w}_{n}^{\prime} \mathbf{Y}_{n i}\right), \quad 1 \leq i \leq n
$$

as a consequence of a change of the variance-covariance matrix $\boldsymbol{\Sigma}_{n}$ in a high-dimensional setting.

## Projection-Based Analysis

To proceed, let us more generally consider the quadratic form

$$
Q_{n}\left(\mathbf{v}_{n}, \mathbf{w}_{n}\right)=\mathbf{v}_{n}^{\prime} \boldsymbol{\Sigma}_{n} \mathbf{w}_{n}
$$

for such weighting vectors $\mathbf{v}_{n}$ and $\mathbf{w}_{n}$.
Remark: Observe that even for $\boldsymbol{\Sigma}_{n}=\sigma \mathbf{1 1}^{\prime}$ we have

$$
\left|Q_{n}\left(\mathbf{v}_{n}, \mathbf{w}_{n}\right)\right|=\sigma\left|\mathbf{v}_{n}^{\prime} \mathbf{1 1 ^ { \prime }} \mathbf{w}_{n}\right|=\sigma\left|\sum_{i} v_{n i} \sum_{i}\right| w_{n i} \mid \leq \sigma\left\|\mathbf{v}_{n}\right\|_{\ell_{1}}\left\|\mathbf{w}_{n}\right\|_{\ell_{1}}
$$

So, the $\ell_{1}$ condition is a natural one and ensures that even full covariance matrices are not mapped to $\infty$.

## Framework and Assumptions

Model: The coordinates are linear processes

$$
Y_{k}^{(\nu)}=Y_{n k}^{(\nu)}=\sum_{j=0}^{\infty} c_{n j}^{(\nu)} \epsilon_{k-j}, \quad k=1, \ldots, n
$$

for coefficients $\left\{c_{n j}^{(\nu)}: j \in \mathbb{N}_{0}\right\}, \nu=1, \ldots, d_{n}$, and mean zero independent r.v.s. $\left\{\epsilon_{k}\right\}$ with

$$
E\left|\epsilon_{k}\right|^{4+\delta}<\infty
$$

for some $\delta>0$.
Assumption A: The sequences $\left\{c_{n j}^{(\nu)}: j \in \mathbb{N}_{0}\right\}$ satisfy

$$
\begin{equation*}
\sup _{n \in \mathbb{N}} \max _{1 \leq \nu \leq d_{n}}\left|c_{n j}^{(\nu)}\right|^{2} \ll j^{-3 / 2-\theta / 2} \tag{3}
\end{equation*}
$$

for some $0<\theta<1 / 2$.

## Strong Approximation

Define

$$
\begin{align*}
& \widehat{\boldsymbol{\Sigma}}_{n k}=\left(\sum_{i=1}^{k} Y_{i}^{(\nu)} Y_{i}^{(\mu)}\right)_{1 \leq \nu, \mu \leq d_{n}},  \tag{4}\\
& \boldsymbol{\Sigma}_{n k}=\left(\sum_{i=1}^{k} E Y_{i}^{(\nu)} Y_{i}^{(\mu)}\right)_{1 \leq \nu, \mu \leq d_{n}}, \tag{5}
\end{align*}
$$

for $n, k \geq 1$. To be precise, our results shall deal with

$$
D_{n k}=\mathbf{v}_{n}^{\prime}\left(\widehat{\boldsymbol{\Sigma}}_{n k}-\boldsymbol{\Sigma}_{n k}\right) \mathbf{w}_{n}, \quad n, k \geq 1
$$

and the associated càdlàg processes

$$
\mathcal{D}_{n}(t)=\mathbf{v}_{n}^{\prime} n^{-1 / 2}\left(\widehat{\boldsymbol{\Sigma}}_{n,\lfloor n t\rfloor}-\boldsymbol{\Sigma}_{n,\lfloor n t\rfloor}\right) \mathbf{w}_{n}, \quad t \in[0,1], n \geq 1 .
$$

If the dependence of the above quantities on $\mathbf{v}_{n}, \mathbf{w}_{n}$ matters, we shall indicate this in our notation and then write

$$
D_{n k}\left(\mathbf{v}_{n}, \mathbf{w}_{n}\right), \mathcal{D}_{n}\left(t ; \mathbf{v}_{n}, \mathbf{w}_{n}\right) .
$$

## Strong Approximation

Recalling that $\widehat{\boldsymbol{\Sigma}}_{n}=n^{-1} \widehat{\boldsymbol{\Sigma}}_{n, n}$, cf. (1) and (4), we have

$$
\mathcal{D}_{n}(1)=\mathbf{v}_{n}^{\prime} \sqrt{n}\left(\widehat{\boldsymbol{\Sigma}}_{n}-\boldsymbol{\Sigma}_{n}\right) \mathbf{w}_{n}, \quad n \geq 1
$$

is the centered and scaled version of the bilinear form

$$
Q\left(\mathbf{v}_{n}, \mathbf{w}_{n}\right)=\mathbf{v}_{n}^{\prime} \widehat{\boldsymbol{\Sigma}}_{n} \mathbf{w}_{n}=\widehat{\operatorname{Cov}}\left(\mathbf{v}_{n}^{\prime} \mathbf{Y}_{n}, \mathbf{w}_{n}^{\prime} \mathbf{Y}_{n}\right)
$$

where

$$
\boldsymbol{\Sigma}_{n}=E \widehat{\boldsymbol{\Sigma}}_{n}=\frac{1}{n} \sum_{i=1}^{n} E\left(\mathbf{Y}_{n i} \mathbf{Y}_{n i}\right)^{\prime}
$$

If $\left\{\mathbf{Y}_{n i}: 1 \leq i \leq n\right\}$ is stationary, then $\boldsymbol{\Sigma}_{n}$ simplifies to $\boldsymbol{\Sigma}_{n}=E\left(\mathbf{Y}_{n 1} \mathbf{Y}_{n 1}^{\prime}\right)$ (but our result are more general).

## Strong Approximation

## Result:

Within the model framework under certain additional technical conditions, we may approximate the related processes by Brownian motions:

$$
\left|D_{n t}-\alpha_{n} B_{n}(t)\right|=o\left(t^{1 / 2}\right), \quad \text { for all } t>0 \text { a.s. }
$$

as $n, t \rightarrow \infty$, and

$$
\sup _{t \in[0,1]}\left|\mathcal{D}_{n}(t)-\alpha_{n} B_{n}(\lfloor n t\rfloor / n)\right|=o(1), \quad \text { a.s. }
$$

as $n \rightarrow \infty$, as well as the CLT

$$
\left|\mathcal{D}_{n}(1)-\alpha_{n} B_{n}(1)\right|=o(1), \quad \text { a.s. }
$$

as $n \rightarrow \infty$, i.e. $\mathcal{D}_{n}(1)$ is asymptotically $\mathcal{N}\left(0, \alpha_{n}^{2}\right)$.

## Strong Approximation

## Standardized sequential statistic:

Monitor the sequence of (standardized) deviations from an assumed variance-covariance matrix $\boldsymbol{\Sigma}_{n}$ via

$$
\mathcal{D}_{n}^{*}(t)=\alpha_{n}^{-1}\left(\mathbf{v}_{n}, \mathbf{w}_{n}\right) \mathcal{D}_{n}\left(t, \mathbf{v}_{n}, \mathbf{w}_{n}\right), \quad t \in\left[s_{0}, 1\right]
$$

which can be approximated by a Brownian motion for large $n$.
Those results provide a basis for valid statistical inference.

## Strong Approximation

## Multivariate Extension:

Needed when projecting high-dimensional data in lower-dimensional subspaces!

## Theorem

Let $\left\{\mathbf{v}_{n j}, \mathbf{w}_{n j}: 1 \leq j \leq K\right\}$ be weighting vectors of dimension $d_{n}$ satisfying condition (7).
Then, under the assumptions of the previous theorem, there exists a $K$-dimensional Brownian motion $\left\{\mathbf{B}^{(n)}(t): t \in[0,1]\right\}$ with coordinates $B_{n i}=B_{n}\left(t ; \mathbf{v}_{n i}, \mathbf{w}_{n i}\right), t \in[0,1], i=1, \ldots, K$, such that

$$
\begin{equation*}
\left\|\left(\mathcal{D}_{n}\left(t ; \mathbf{v}_{n i}, \mathbf{w}_{n i}\right)\right)_{i=1}^{K}-\left(B_{n}\left(\lfloor n t\rfloor / n ; \mathbf{v}_{n i}, \mathbf{w}_{n i}\right)\right)_{i=1}^{K}\right\|=o(1), \tag{6}
\end{equation*}
$$

a.s., as $n \rightarrow \infty$, where $\|\bullet\|$ denotes an arbitrary vector norm on $\mathbb{R}^{K}$.

## Strong Approximation

## Corollary

Suppose that $\mathbf{Y}_{n 1}, \ldots, \mathbf{Y}_{n n}$ is a $d_{n}$-dimensional vector time series satisfying Assumption (A). Then, after redefining the series on a new probability space, there exists a Brownian motion such that

$$
\left|\max _{k \leq n}\right| \mathcal{D}_{n}(k / n)\left|-\max _{k \leq n}\right| \alpha_{n} B_{n}(k / n)| |=o(1)
$$

as $n \rightarrow \infty$.
The proofs rely on generalizations of Kouritzin (1995, SPA), who applied Philipp's (1986) results on strong approximations in Hilbert spaces.

## Discussion

...of the $\ell_{1}$-condition:

$$
\begin{equation*}
\sup _{n \in \mathbb{N}}\left\|\mathbf{w}_{n}\right\|_{\ell_{1}}=\sup _{n \in \mathbb{N}} \sum_{\nu=1}^{d_{n}}\left|w_{j}\right|<\infty \tag{7}
\end{equation*}
$$

Ex. 1: $\ell_{0}$-sparse vectors: $w_{i}>0$ only for $i \in\left\{i_{1}, \ldots, i_{L}\right\}, L$ fixed. (classical 'low-dimensional' case)
Ex. 2: $w_{n}^{\prime}=\left(w_{1}, \ldots, w_{d_{n}}\right)^{\prime}$ with $\sum_{j}$ (most coordinates receive a negligible weight)
Ex. 3: $w_{n i}=1 / d_{n}$ for $i=1, \ldots, d_{n}$.
(all $d_{n}$ coordinates are taken into account.)

## Discussion

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(all $d_{n}$ coordinates are taken into account.)

## Discussion

... of Assumption A:

$$
\begin{equation*}
\sup _{n \in \mathbb{N}} \max _{1 \leq \nu \leq d_{n}}\left|c_{n j}^{(\nu)}\right|^{2} \ll j^{-3 / 2-\theta / 2} \tag{8}
\end{equation*}
$$

for some $0<\theta<1 / 2$.
Assumption A...

- $c_{n j}^{(\nu)}=c_{j}^{(\nu)}$ : At time $n$ we observe $d_{n}$ sequences (not depending on $n$ ). But we allow for arrays.
- covers various short memory processes, e.g. ARMA processes.
- covers many long-range dependent series such as fractionally integrated noise of order $d \in(-1 / 2,1 / 4-\theta / 2)$,

$$
(1-L)^{d} X_{t}=\epsilon_{t}
$$

## Discussion

Define the scaled Frobenius norm by

$$
\|\mathbf{A}\|_{F}^{*}=\frac{1}{d_{n}^{1 / 2}}\left(\sum_{i, j=1}^{d_{n}} a_{i j}^{2}\right)^{1 / 2} \quad\left(\text { s.th. }\left\|\mathbf{I}_{d_{n}}\right\|_{F}^{*}=1\right)
$$

## Lemma

Suppose Assumption (A) holds true and that, for fixed $t$, the variances $\sigma_{t-j}^{2}, j \geq 0$, of the innovations satisfy

$$
\sum_{j=0}^{\infty} j^{-3 / 2-\theta / 2} \sigma_{t-j}^{2}<\infty
$$

Then, as $r \rightarrow \infty$; we have

$$
\sup _{n \in \mathbb{N}}\left\|\boldsymbol{\Sigma}_{n}[t]-\sum_{j=1}^{r} \sigma_{t-j}^{2} \mathbf{c}_{n j} \mathbf{c}_{n j}\right\|^{\prime} \|_{F}^{*}=o(1) .
$$

## Applications: Portfolio Selection

Consider Assets returns $\mathbf{R}_{n}=\left(R_{n}^{(1)}, \ldots, R_{n}^{\left(d_{n}\right)}\right)^{\prime}$ corresponding to the time period $[n-1, n]$ with mean vector $\boldsymbol{\mu}$ and covariance matrix $\boldsymbol{\Sigma}=\left(\sigma_{i j}\right)_{i j}$.
Since $\sigma_{i j}$ is the covariance between the return of asset $i$ and asset $j, 1 \leq i, j, \leq d_{n}$, it is not restrictive to assume that the entries of $\boldsymbol{\Sigma}$ neither depend on $n$ nor $d_{n}$.
An investor holds at time $n-1$ the position $w_{n j}$ in asset $j$. $w_{n j}>0$ long position, $w_{n j}<0$ short position.
W.I.o.g. the initial value (capital) at time $n-1$ equals $V=\sum_{j=1}^{d_{n}} w_{n j}=1$.
Then the value at time instant $n$ is $\mathbf{w}_{n}^{\prime} \mathbf{R}_{n}$.

## Applications: Portfolio Selection

Classical formulation of the portfolio selection problem: Risk $=$ Variance:

$$
\min _{\mathbf{w}_{n}} \operatorname{Var}\left(\mathbf{w}_{n}^{\prime} \mathbf{R}_{n}\right)=\mathbf{w}_{n}^{\prime} \boldsymbol{\Sigma} \mathbf{w}_{n}, \quad \text { subject to } \mathbf{w}_{n}^{\prime} \mathbf{1}=1
$$

whose solution is known to be

$$
\mathbf{w}_{n}^{* \prime}=\left(\mathbf{1}^{\prime} \Sigma^{-1} \mathbf{1}\right)^{-1} \mathbf{1}^{\prime} \Sigma^{-1}
$$

If that solution satisfies the no-short-sales condition, then $\left\|\mathbf{w}_{n}^{*}\right\|_{\ell_{1}}=\mathbf{1}^{\prime} \mathbf{w}_{n}^{*}=1$.
Provided the vector time series of returns satisfies our assumptions, our results provide the asymptotics for the optimal risk

$$
\operatorname{Var}\left(\left(\mathbf{w}_{n}^{*}\right)^{\prime} \mathbf{Y}_{n}\right)
$$

 $\mathbf{w}_{n}^{*}$ ) from an independent learning sample.

## Applications: Shrinkage

Shrinkage estimation is a well established approach for regularization (Ledoit \& Wolf (2004); Sancetta, 2008).
To improve properties such as $E\left\|\widehat{\boldsymbol{\Sigma}}_{n}-\boldsymbol{\Sigma}_{n}\right\|_{F}^{2}$ or the condition number, one estimates $\boldsymbol{\Sigma}_{n}$ by a linear (convex) combination of $\widehat{\boldsymbol{\Sigma}}_{n}$ and a well-conditioned target such as the identity.
Projecting $\boldsymbol{\Sigma}_{n}$ onto span $\left\{\mathrm{id}_{n}\right\}$ leads to the target $\boldsymbol{\Sigma}_{n}^{(0)}=\mu_{n} \mathrm{id}_{n}$, where $\mu_{n}=\operatorname{tr}\left(\boldsymbol{\Sigma}_{n}\right)$ (shrinkage intensity).

We are led to the shrinkage estimator

$$
\boldsymbol{\Sigma}_{n}^{s}\left(W_{n}\right)=\left(1-W_{n}\right) \widehat{\boldsymbol{\Sigma}}_{n}+W_{n} \mu_{n} \operatorname{id}_{n}
$$

Optimizing $W_{n}$ w.r.t. the MSE

$$
W_{n}^{*}=\operatorname{argmin}_{W_{n} \in[0,1]} d_{n}^{-1} E\left\|\boldsymbol{\Sigma}_{n}^{s}\left(W_{n}\right)-\boldsymbol{\Sigma}_{n}\right\|_{F}^{2}
$$

leads to explitic formulas for $W_{n}^{*}$ and ensures a true improvement

$$
\begin{equation*}
E\left\|\boldsymbol{\Sigma}_{n}^{s}-\boldsymbol{\Sigma}_{n}\right\|_{F}^{2}<E\left\|\widehat{\boldsymbol{\Sigma}}_{n}-\boldsymbol{\Sigma}_{n}\right\|_{F}^{2} \tag{ISW}
\end{equation*}
$$

## Applications: Sparse Principal Component Analysis LASSO

Let $\mathcal{X}_{n}$ be a $\left(n \times d_{n}\right)$-dimensional data matrix, independent from $\mathcal{Y}_{n}$.

SCotLASS (Simplified component technique-lasso), Jolliffee (2003): $1^{\text {st }}$ principal component (pc) solves

$$
\max _{\mathbf{v}} \mathbf{v}^{\prime} \mathcal{X}_{n}^{\prime} \mathcal{X}_{n} \mathbf{v}, \quad \text { subject to }\|\mathbf{v}\|_{\ell_{2}}^{2} \leq 1,\|\mathbf{v}\|_{\ell_{1}} \leq c
$$

Continue in this way under the additional constraints that further components are orthogonal.

## Applications: Sparse Principal Component Analysis

LASSO (Tibshirani, 1996 \& 2011): Determine $\ell_{1}$-sparse coefficient vector in a high-dimensional linear regression in dim. $p_{n}$

$$
Y_{t}=\mathbf{X}_{t}^{\prime} \beta_{0}+\epsilon_{t}, \quad E\left(\epsilon_{t} \mid \mathbf{X}_{t}\right)=0
$$

Given an estimator $\widehat{\beta}_{n}, \pi_{n}=\mathbf{X}^{\prime} \widehat{\beta}_{n}$ is used for prediction.
LASSO minimizes the $\ell_{1}$-constrained least squares criterion

$$
\beta \mapsto \sum_{t}\left(Y_{t}-\mathbf{X}_{t}^{\prime} \beta\right)^{2},\|\beta\|_{\ell_{1}} \leq c
$$

for some bound $c>0$.
Apply results with $\mathbf{w}_{n}=\widehat{\beta}_{n}$ estimated from indep. learning sample, $d_{n}=p_{n}, \mathbf{Y}=\mathbf{X}, \mathbf{Y}_{t}=\mathbf{X}_{t}$ to infer $\operatorname{Var}\left(\mathbf{w}_{n}^{\prime} \mathbf{X}\right)$ given the learning sample, provided $\left\{\mathbf{X}_{t}\right\}$ satisfies our assumptions.

## End

## Related Publication:

Steland, A. and R. v. Sachs (2015). Large sample approximations for variance-covariance matrices of high-dimensional time series, under revision.

Thanks for your attention.

