An Algebraic Estimator for Large Spectral Density Matrices

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\textbf{ABSTRACT}

We propose a new estimator of high-dimensional spectral density matrices, called ALgebraic Spectral Estimator (ALSE), under the assumption of an underlying low rank plus sparse structure, as typically assumed in dynamic factor models. The ALSE is computed by minimizing a quadratic loss under a nuclear norm plus $l_1$ norm constraint to control the latent rank and the residual sparsity pattern. The loss function requires as input the classical smoothed periodogram estimator and two threshold parameters, the choice of which is thoroughly discussed. We prove consistency of ALSE as both the dimension $p$ and the sample size $T$ diverge to infinity, as well as the recovery of latent rank and residual sparsity pattern with probability one. We then propose the UNshrunk ALgebraic Spectral Estimator (UNALSE), which is designed to minimize the Frobenius loss with respect to the pre-estimator while retaining the optimality of the ALSE. When applying UNALSE to a standard U.S. quarterly macroeconomic dataset, we find evidence of two main sources of comovements: a real factor driving the economy at business cycle frequencies, and a nominal factor driving the higher frequency dynamics. The article is also complemented by an extensive simulation exercise. Supplementary materials for this article are available online.

\textbf{1. Introduction}

An appealing, natural, and classical way to model time series data is through spectral analysis (Brillinger 2001). Given a $p$-dimensional vector stochastic process, its $p \times p$ spectral density matrix characterizes all second order dependencies. Moreover, conditional second order dependencies can also be extracted starting from the inverse of the spectral density matrix. The spectral approach is appealing since, once we move from the time domain to the frequency domain, data become asymptotically independent, as the sample size $T$ grows to infinity.

Statistical methods for the study of time series based on spectral analysis include: pseudo-maximum likelihood estimation (Velasco and Robinson 2000), linear regression (Harvey 1978), cointegration tests or information criteria based on the zero-frequency spectral density matrix of a vector of time series (Stock and Watson 1988), and similarly seasonal cointegration tests based on the spectral density matrix at selected frequencies (Joyeux 1992), de-trending methods (Corbæ, Ouliaris, and Phillips 2002), Granger causality tests (Breitung and Candelon 2006; Farnè and Montanari 2021), and the analysis of low frequency co-movements (Müller and Watson 2018). Finally, the inverse spectral density matrix is at the basis of graphical models and dynamic network analysis (Dahlhaus 2000; Eichler 2007; Davis, Zang, and Zheng 2016).

The use of spectral analysis is widespread in many applied fields. Examples are the construction of business cycle indicators in macroeconomics (Altissimo et al. 2010), portfolio optimization at different horizons in finance (Chaudhuri and Lo 2015), and the study of brain activity in biostatistics (Ombao et al. 2001; Ombao, Von Sachs, and Guo 2005; Fiecas and Ombao 2011, 2016).

All above methods and applications require as input an estimator of the spectral density matrix or of its inverse. Just like for the covariance matrix estimation in time domain, estimation of a spectral density matrix is a hard problem when the dimension of the process $p$ is comparable, or even larger, than the sample size $T$. In this case, the classical smoothed periodogram estimator is not positive definite simply due to lack of degrees of freedom. Given the increased availability of large datasets in the recent years, this issue becomes of fundamental importance. Wu and Zaffaroni (2018) provide consistency results for the smoothed periodogram estimator in high dimension, which hold uniformly over all frequencies. Zhang and Wu (2021) improve on those results, by providing bounds for the estimation error which hold uniformly over all $p^2$ entries.

To solve the problem of the curse of dimensionality, here, we start from observing that the second moments of most high-dimensional time series tend to have both a low rank and a sparse component. Indeed, on the one hand, most economic datasets are known to be mainly “dense” rather than sparse (Giannone, Lenza, and Primiceri 2021). Moreover, there exist mathematical results proving that large dimensional panels of time series can in general be represented as having a factor structure (Forni and Lippi 2001; Hallin and Lippi 2013). On the other hand, once the common factors are controlled for, there is evidence of sparseness in the second order structure of the residuals (Barigozzi and Hallin 2017).
In this article, we assume that the spectral density matrix, \( \Sigma(\theta) \), of \( p \)-dimensional time series has the low rank plus sparse structure: \( \Sigma(\theta) = L(\theta) + \delta S(\theta), \theta \in [-\pi, \pi] \), where \( L(\theta) \) has rank \( r < p \), and \( S(\theta) \) is a sparse matrix. Based on this assumption, our estimators \( \hat{L}(\theta) \) and \( \hat{S}(\theta) \) of the two components of the spectral density matrix are obtained by regularizing the smoothed periodogram estimator, \( \Sigma(\theta) \), by means of a nuclear norm plus \( l_1 \) norm penalization. Specifically, at each given frequency \( \theta \in [-\pi, \pi] \) we define the ALgebraic Spectral Estimator (ALSE) as the couple of matrices

\[
(\hat{L}(\theta), \hat{S}(\theta)) = \arg\min_{L(\theta), S(\theta)} \frac{1}{2} \| \Sigma(\theta) - (L(\theta) + S(\theta)) \|_F^2 + \psi \| L(\theta) \|_* + \rho \| S(\theta) \|_1,
\]

where \( L(\theta) \) and \( S(\theta) \) indicate generic values of the matrices belonging to appropriate algebraic matrix varieties, \( \| L(\theta) \|_* = \text{tr}(L(\theta)) \) and \( \| S(\theta) \|_1 = \sum_{i,j} | S_{ij}(\theta) | \), and \( \psi \) and \( \rho \) are threshold parameters. An estimator of the spectral density matrix is then \( \Sigma(\theta) = \hat{L}(\theta) + \hat{S}(\theta) \). The above optimization problem is solved by iterating between a singular value thresholding step (Daubechies, Defrise, and DeMol 2004), giving \( \hat{S}(\theta) \), and then, in a second step, apply soft or hard thresholding to the orthogonal complement to obtain a sparse and positive definite residual. Their resulting estimator is called POET. Farnè and Montanari (2020) adopt a minimization algorithm analogous to the one considered in this article which recovers the covariance matrix consistently, both algebraically and parametrically. The resulting estimator is called UNALCE and they show that it systematically outperforms POET both in terms of parametric consistency, and, more importantly, because it provides the algebraic recovery of latent rank and sparsity pattern. A similar approach was proposed by Luo (2011), however, it is based on the assumption of bounded eigenvalues for the covariance matrix, which does not allow for the joint identification of the two components.

We then apply to ALSE in (1) the unshrinkage step of estimated latent eigenvalues by Farnè and Montanari (2020), which optimizes the finite sample Frobenius loss with respect to the smoothed periodogram while retaining algebraic consistency. This step consists in unshrinkifying \( \hat{L}(\theta) \), while a new estimate of the sparse component is obtained by retaining the off-diagonal sparsity pattern of \( \hat{S}(\theta) \), and computing its diagonal by difference from the diagonal of \( \hat{S}(\theta) \) and the new estimate of the low rank component. We call the resulting estimator of the overall spectral density matrix UNShrunken ALgebraic Spectral Estimator (UNALSE). By construction UNALSE improves over ALSE in terms of Frobenius norm while it is equivalent to ALSE in terms of spectral norm.

Our approach is based on the fundamental identification assumptions we make on the behavior of the eigenvalues of the spectral density matrix. We assume the \( r \) eigenvalues of the low rank component, \( L(\theta) \), to be diverging at a rate \( p^\alpha \) with \( \alpha \in [0, \frac{1}{2}] \). In the language of factor models, this means we are allowing for the presence of factors with different degrees of pervasiveness, that is, both weak and strong factors. Moreover, we assume the sparse component, \( S(\theta) \), to have eigenvalues diverging at most at a rate \( p^\delta \) with \( \delta \in [0, 1/2] \) and \( \delta < \alpha \). These assumptions imply the existence of an eigen-gap in the spectrum of the spectral density matrix, \( \Sigma(\theta) \), which widens as \( p \) increases.

There exist alternative approaches to the estimation of large spectral density matrices. Forni et al. (2000) propose principal component analysis in the frequency domain to recover the low rank component. Böhm and von Sachs (2008, 2009) propose to shrink the smoothed periodogram toward either a reduced rank target or the identity, respectively. Fiecas and von Sachs (2014) propose a penalized likelihood approach, and Fiecas et al. (2019) consider constrained \( l_1 \) minimization for estimating the inverse. While some of those works assume either a low rank or a sparsity structure, none of them considers both assumptions jointly.

Similar approaches based on a low rank plus sparse assumption exist in time domain, that is, for the estimation of the covariance matrix. Fan, Liao, and Mincheva (2013) consider principal components to recover the low rank component and then, in a second step, apply soft or hard thresholding to the orthogonal complement to obtain a sparse and positive definite residual. Their resulting estimator is called POET. Farnè and Montanari (2020) adopt a minimization algorithm analogous to the one considered in this article which recovers the covariance matrix consistently, both algebraically and parametrically. The resulting estimator is called UNALCE and they show that it systematically outperforms POET both in terms of parametric consistency, and, more importantly, because it provides the algebraic recovery of latent rank and sparsity pattern. A similar approach was proposed by Luo (2011), however, it is based on the assumption of bounded eigenvalues for the covariance matrix, which does not allow for the joint identification of the two components.

Our assumption of a low rank plus sparse decomposition of the spectral density matrix is also strictly related to, and inspired by, the Generalized Dynamic Factor Model (GDFM) representation of a large panel of time series, originally proved by Forni and Lippi (2001). This is a very popular approach to dimension reduction (see, e.g., the application in Altissimo et al. 2010). In the GDFM, \( r \) latent factors are loaded by each series in a dynamic way, that is, not only contemporaneously but also with lags. The key assumptions are: (i) pervasiveness of the factors resulting in \( r \) leading spiking spectral eigenvalues, and (ii) weak serial and cross-correlation in the residuals, resulting in boundedness of the spectral eigenvalues. In our notation these conditions imply \( \alpha = 1 \) and \( \delta = 0 \).

Forni et al. (2000, 2005, 2017) consider different estimators of the GDFM, which are all built starting from a consistent estimator of the spectral density matrix. In particular, in all those approaches the low rank component of the spectral density is estimated via the \( r \) leading dynamic principal components, that is, the principal components of the spectral density matrix across frequencies of the smoothed periodogram (see also Brillinger 2001). The consistency of this method relies on the pervasiveness of spectral eigenvalues with respect to the dimension \( p \). The spectral density of the residual component, called idiosyncratic component in the GDFM literature, is then estimated as the difference between the estimated spectral density of the observed data and its estimated low rank component. Hence, by construction, the spectral density of the idiosyncratic
component has rank $p - r$, that is, it is not positive definite, and, therefore, not invertible. There exist also few papers dealing with determining the dynamic rank, $r$: Hallin and Liška (2007) propose an information criterion, and Onatski (2009) proposes a test based on the asymptotic distribution of the spectral eigenvalues.

The above approaches to the estimation of the GDFM suffer from some drawbacks. First, any estimator of the spectral density matrix based on the principal components of an input estimator, like the smoothed periodogram, is likely to suffer from numerical instability, especially if $p$ is large, due to the Marcenko and Pastur (1967) law. Second, the strict perversiveness assumption of spectral eigenvalues ($\omega = 1$) is rarely satisfied in practice, since the factor strength might vary across frequencies, for example, due to common, frequency specific, features. Third, the weak correlation assumption increases the number of parameters when $p$ is large, which prevents the residual component to be identified.

The estimator we propose in this article is able to address those drawbacks, because it tolerates weakly pervasive factors, relevantly reduces the number of estimated parameters and, given its algebraic consistency, it is also a consistent estimator of the latent rank $r$. For these reasons, it can be used as input of all the estimators of the GDFM considered in the literature.

The article is organized as follows. In Section 2 we present our main results using the GDFM setting as a guiding example. In Sections 3, 4, and 5 we present the general framework, describe estimation, and prove consistency. In Section 6 we discuss the unshrinking of ALSE generating UNALSE. In Section 7 we discuss the choice of the threshold parameters. Section 8 presents the results when applying UNALSE to a dataset of quarterly U.S. macroeconomic time series. In Appendix A, supplementary materials we prove all theoretical results of the article. In Appendix B, supplementary materials we consider the implications of our assumptions for a large VAR setting. In Appendix C, supplementary materials we show how the theory presented can be extended also to nonlinear models. In Appendix D, supplementary materials we show simulation results under a variety of data generating processes.

Notation. We denote a $p \times p$ Hermitian positive-definite complex matrix as $A$ and its transposed complex conjugate as $A^\dagger$. Let $\lambda_j(A)$, $i = 1, \ldots, p$, be the (real) eigenvalues of $A$ in descending order, and by $A_{ij}$ its $(i,j)$th entry. Define $\overline{A}_{ij}$ as the complex conjugate of $A_{ij}$, and $|A_{ij}| = \sqrt{A_{ij}\overline{A}_{ij}}$, the real and imaginary parts are indicated as $\text{Re}(A_{ij})$ and $\text{Im}(A_{ij})$, respectively. To indicate that $A$ is positive definite or semidefinite we write: $A \succ 0$ or $A \succeq 0$, respectively. We use the following norms. Element-wise norms: $l_0$ norm: $\|A\|_0 = \sum_{i=1}^{p} \sum_{j=1}^{p} 1(A_{ij} \neq 0)$; $l_1$ norm: $\|A\|_1 = \sum_{j=1}^{p} \sum_{i=1}^{p} |A_{ij}|$; Frobenius norm: $\|A\|_F = \sqrt{\sum_{i=1}^{p} \sum_{j=1}^{p} |A_{ij}|^2} = \sqrt{\text{tr}(AA^\dagger)}$; maximum norm: $\|A\|_\infty = \max_{1 \leq i,j \leq p} |A_{ij}|$. Vector-induced norms: $l_0,v$ norm $\|A\|_{0,v} = \max_{1 \leq i \leq p} \sum_{j=1}^{p} 1(A_{ij} \neq 0)$; $l_1,v$ norm $\|A\|_{1,v} = \max_{1 \leq i \leq p} \sum_{j=1}^{p} |A_{ij}|$; $l_{\infty,v}$ norm $\|A\|_{\infty,v} = \max_{1 \leq i \leq p} \sum_{j=1}^{p} |A_{ij}|$; spectral norm: $\|A\|_2 = \sqrt{\lambda_1(AA^\dagger)} = \lambda_1(A)$; nuclear norm: $\|A\|_* = \text{tr}(A) = \sum_{i=1}^{p} \lambda_i(A)$.

For two sequences $\{a_n\}$ and $\{b_n\}$ such that $a_n, b_n \to \infty$ as $n \to \infty$, we write $a_n = O(b_n)$ if $a_n/b_n \leq C$ for some finite positive real $C$ independent of $n$ and $a_n \asymp b_n$ if $a_n/b_n = O(1)$ and $b_n/a_n = O(1)$.

2. Overview of Results

In this section, we present the main features of our estimator under the assumption that the data follow a Generalized Dynamic Factor Model (GDFM) as defined by Forni and Lippi (2001). The GDFM setting has to be considered just as a motivating example, which is well suited to allow the reader to immediately appreciate the contribution of this article with respect to the state of art. In the following sections, we present our theory in more detail showing that the validity of our results is actually much broader than the case here considered.

Let $X = \{X_t, i = 1, \ldots, p, t \in \mathbb{Z}\}$ be a $p$-dimensional panel of zero-mean second-order stationary time series. The set of all $L_2$-convergent linear combinations of $X_t$'s and their limits, as $p \to \infty$, of $L_2$-convergent sequences thereof, is a Hilbert space, denoted by $H_X$. Hence, for all $t \in \mathbb{Z}$ and all $p \in \mathbb{N}$, any dynamic linear combination of $X_t$'s, $y_t = \sum_{i=1}^{p} \sum_{k=-\infty}^{k=\infty} a_{ik}X_{t-i-k}$, such that $\sum_{i=1}^{p} \sum_{k=-\infty}^{k=\infty} a_{ik}^2 = 1$, belongs to $H_X$. Following Hallin and Lippi (2013, Definitions 2.1 and 2.2), we define as common variable the $L_2$-limit of any standardized dynamic linear combination of the $X_t$'s, say $\frac{y_t}{\sqrt{\text{var}(y_t)}}$, such that $\text{var}(y_t) \to \infty$, as $p \to \infty$. The Hilbert space of all common variables is denoted by $H_{\text{com}}$, while its orthogonal complement with respect to $H_X$, denoted as $H_{\text{idio}}$, contains all the idiosyncratic variables, that is, all dynamic linear combinations $y_t$ with bounded variance $\text{var}(y_t)$ for all $p \in \mathbb{N}$.

Hallin and Lippi (2013) prove that there exist two unique zero-mean stochastic processes $\{\chi_t\} \in H_{\text{com}}$ and $\{\epsilon_t\} \in H_{\text{idio}}$, mutually orthogonal at all leads and lags, such that

$$X_t = \chi_t + \epsilon_t, \quad i \in \mathbb{N}, \quad t \in \mathbb{Z}.$$  (2)

The process $\{\chi_t\}$ is called common component, the process $\{\epsilon_t\}$ is called idiosyncratic component. We denote the $p \times p$ spectral density matrices of the $p$-dimensional processes $\{X_t = (X_{1t} \ldots X_{pt})\}$, $\{\chi_t = (\chi_{1t} \ldots \chi_{pt})\}$, and $\{\epsilon_t = (\epsilon_{1t} \ldots \epsilon_{pt})\}$ by $\Sigma(\theta), L(\theta)$ and $S(\theta)$, respectively. Under representation (2): (i) $\{\chi_t\}$ is driven by an $r$-tuple of mutually orthogonal white noises loaded by a linear time filter for all $p \in \mathbb{N}$, that is, $r_k(L(\theta)) = r$ for all $\theta \in [-\pi, \pi]$; (ii) $\{\epsilon_t\}$ is orthogonal to those white noises at all leads and lags, and (iii) $\{\chi_t\}$ follows representation (2) if and only if the $r$ eigenvalues of $L(\theta)$ diverge for all $\theta \in [-\pi, \pi]$ as $p$ diverges, while the eigenvalues of $S(\theta)$ remain bounded for all $p \in \mathbb{N}$ and all $\theta \in [-\pi, \pi]$ (Forni et al. 2000). The scalar $r$ is called the number of dynamic factors. All this defines the GDFM, which encompasses the approximate static factor models of Chamberlain and Rothschild (1983), as well as the exact dynamic factor models of Sargent and Sims (1977).

As usual in the GDFM literature, in this section we adopt the assumption (relaxed later on) that the $r$ eigenvalues of $\frac{L(\theta)}{p}$ are finite and bounded away from zero for all $p \in \mathbb{N}$ and all $\theta \in (-\pi, \pi)$.
Similarly, the definition of idiosyncratic variable leads to the condition \( \lambda_1(S(\theta)) = \| S(\theta) \|_2 < \infty \) for all \( p \in \mathbb{N} \) and all \( \theta \in [-\pi, \pi] \). These assumptions on \( L(\theta) \) and \( S(\theta) \) imply that the gap between the \( r \)-th and the \( (r+1) \)-th eigenvalue of the spectral density matrix \( \Sigma(\theta) \) increases at all \( \theta \in [-\pi, \pi] \) as \( p \) increases, making the recovery of the low rank component possible in the limit \( p \to \infty \).

In this section, we further adopt the assumption (relaxed later on) that the idiosyncratic spectral density matrix \( S(\theta) \) is such that \( \| S(\theta) \|_{0,v} \) is bounded for all \( p \in \mathbb{N} \) and all \( \theta \in [-\pi, \pi] \). Since \( \| S(\theta) \|_2 \leq \| S(\theta) \|_{0,v} \), the original assumption \( \| S(\theta) \|_2 < \infty \) still holds. This is done in order to enforce element-wise sparsity on \( S(\theta) \).

Suppose now that we observe a sample of size \( T \) of \( p \)-dimensional data vectors. A classical estimator of the spectral density matrix, which is our pre-estimator, is the smoothed periodogram, defined as

\[
\hat{\Sigma}(\theta_h) = \frac{1}{2\pi} \sum_{k=(T-1)}^{T-1} \left( \frac{k}{MT} \right) e^{-\psi_{th} k} \hat{\Gamma}_X(k),
\]

where \( \hat{\Gamma}_X(k) = T^{-1} \sum_{t=1}^{T-|k|} X_t X_t^T \) and \( K(\cdot) \) is a suitable kernel function with \( MT \) being the associated smoothing parameter. According to Brillinger (2001), for any given element of \( \theta_h \), \( \hat{\Sigma}(\theta_h) \) is consistent if \( \frac{MT}{T} \to 0 \) while \( MT \to \infty \) and \( T \to \infty \). Wu and Zaffaroni (2018) prove the element-wise consistency of \( \hat{\Sigma}(\theta_h) \) uniformly over the frequencies, under appropriate assumptions to be discussed later.

Under the GDFM setting described above, augmented with the sparsity assumption for \( S(\theta) \), we define the ALgebraic Spectral Estimator (ALSE) estimator of the spectral density matrix \( \Sigma(\theta) \) as \( \hat{\Sigma}(\theta) = \hat{L}(\theta) + \hat{S}(\theta) \), where \( \hat{L}(\theta) \) and \( \hat{S}(\theta) \) are such that they satisfy (1). We refer to Section 4 for details on computing the solution of such minimization problem.

Consistency of ALSE under the GDFM setting is in the following Corollary to our main contribution which is Theorem 5.1.

**Corollary 2.1.** For all \( p \in \mathbb{N} \), assume that: (i) the \( r \) nonzero eigenvalues of \( L(\theta) \) are such that \( \frac{\lambda_{r}(L(\theta))}{p} \) is finite and bounded away from zero for all \( j = 1, \ldots, r \); (ii) \( \| S(\theta) \|_{0,v} \) is bounded; (iii) \( \{X_t, t \in \mathbb{Z} \} \) and \( \{\epsilon_t, t \in \mathbb{Z} \} \) are Gaussian processes. Then, under the assumptions and the conditions of Theorem 5.1, there exist real positive \( \kappa_1 \) and \( \kappa_2 \) independent of \( p \) and \( T \) such that, as \( T \to \infty \) and for all \( p \in \mathbb{N} \), the following hold:

\[
1. \quad \mathcal{P} \left( \max_{|h| \leq MT} \frac{1}{p} \| \hat{L}(\theta_h) - L(\theta_h) \|_2 \leq \kappa_1 \sqrt{\frac{MT \log(MT)}{T}} \right) \to 1;
\]

\[
2. \quad \mathcal{P} \left( \| \hat{S}(\theta_h) - S(\theta_h) \|_2 \leq \kappa_2 \left\{ \max_{|h| \leq MT} \| S(\theta_h) \|_{0,v} \right\} \sqrt{\frac{MT \log(MT)}{T}} \right) \to 1;
\]

\[
3. \quad \mathcal{P} \left( \max_{|h| \leq MT} \frac{1}{p} \| \hat{\Sigma}(\theta_h) - \Sigma(\theta_h) \|_2 \leq \kappa_1 \sqrt{\frac{MT \log(MT)}{T}} \right) \to 1.
\]

Furthermore, under the assumptions and the conditions of Corollary 5.1, there exists a positive real \( k_4 \) independent of \( p \) and \( T \) such that, as \( T \to \infty \) and for all \( p \in \mathbb{N} \), the following hold:

\[
\mathcal{P} \left( \max_{|h| \leq MT} \| \hat{S}(\theta_h) - S(\theta_h) \|_2 \leq k_4 \left\{ \max_{|h| \leq MT} \| S(\theta_h) \|_{0,v} \right\} \sqrt{\frac{MT \log(MT)}{T}} \right) \to 1.
\]

The results of Corollary 2.1 contribute to the GDFM literature in three ways. First, the exact rank recovery in part 2 allows to bypass the use of existing criteria for determining the number of factors, like those by Hallin and Liška (2007) and Onatski (2009). Second, we derive a consistency result also for the estimator of the idiosyncratic spectral density \( \hat{S}(\theta) \). Third, assuming that \( S(\theta) \) is full rank, we obtain a result also for the estimators of its inverse. Our consistency rates are comparable to those derived by Fan, Liao, and Mincheva (2013) for the estimation of the covariance matrix of \( \{X_t\} \) when this is generated by a static factor model with sparse idiosyncratic covariance.

In fact Theorem 5.1, which is our main result, holds beyond the standard GDFM assumptions. In particular, first, we relax the strict pervasiveness assumption on latent dynamic factors, by allowing the \( r \) eigenvalues of the matrix \( \frac{L(\theta)}{p} \), with \( \alpha < 1 \), to be finite and bounded away from zero for all \( p \in \mathbb{N} \) and all \( \theta \in [-\pi, \pi] \). Second, we allow for the maximum number of nonzeros per column in \( S(\theta) \), \( \| S(\theta) \|_{0,v} \), to be at most proportional to \( p^\delta \), with \( \delta \in [0, \frac{1}{4}] \) and \( \delta < \alpha \) for all \( p \in \mathbb{N} \) and all \( \theta \in [-\pi, \pi] \). This means that we allow the idiosyncratic spectrum to be quite far from the diagonal matrix. Our setting reduces to the GDFM one when \( \alpha = 1 \) and \( \delta = 0 \). Third, we propose an unshrinking procedure such that the resulting UNShrunken ALgebraic Spectral Estimator (UNALSE) optimizes the finite sample Frobenius loss with respect to the smoothed periodogram while retaining algebraic consistency.

## 3. Model Setup

The aim of this article is estimating the spectral density matrix of a \( p \)-dimensional process \( X = \{X_t, t = 1, \ldots, p \in \mathbb{Z}\} \), following the data generating process:

\[
X_t = \chi_t + \epsilon_t, \quad t \in \mathbb{Z},
\]

\[
\chi_t = \sum_{s=0}^\infty \mathbf{B}_s u_{t-s}, \quad t \in \mathbb{Z},
\]

\[
\epsilon_t = \sum_{s=0}^\infty \mathbf{C}_s \epsilon_{t-s}, \quad t \in \mathbb{Z},
\]

where \( X_t, \chi_t, \epsilon_t, \) and \( \epsilon_t \) are \( p \)-dimensional, \( u_t \) is \( r \)-dimensional, the \( \mathbf{B}_s \) are \( p \times r \), and the \( \mathbf{C}_s \) are \( p \times p \). While Forni and Lippi (2001) derive a two-sided singular MA representation for the process \( \chi_t \), the existence of a one-sided representation (5) is proved by Hallin and Lippi (2013, Theorem 2.2). The MA representation (6) for the process \( \epsilon_t \) is the usual Wold representation.

We make the following assumptions on the MA processes in (5) and (6).

**Assumption 3.1.** (i) \( \{u_t, t \in \mathbb{Z}\} \) is a \( r \)-dimensional independent and identically distributed process with \( \mathbb{E}[u_t] = 0 \), and
E[uiu'i] = I, and with r finite and independent of p for all p ∈ N; (ii) there exists K' > 0 and d_4 > 4 independent of j and t such that E[|uij|^4] ≤ K_a for all j = 1, . . . , r; (iii) for all p ∈ N, (e_t, t ∈ Z) is a p-dimensional independent and identically distributed process with E[e_t] = 0, and E[e_t'e_t] = I_p; (iv) there exists K_e > 0 and d_e > 4 independent of j and t such that E[|e_t|^4] ≤ K_e for all t ∈ N; (v) {u_t} and {e_t} are two mutually independent processes.

Assumption 3.2. There exist M_x, M_e > 0, ρ_x, ρ_e ∈ (0, 1], α ∈ (1/2, 1) and δ' ∈ (0, α), such that, for all p ∈ N: (i) \|B_1\|_∞,∞ ≤ M_x ρ^{\delta'}/2; (ii) \|B_1\|_2 ≤ M_x ρ^{\delta'}/2; (iii) \|C_s\|_{∞,∞} ≤ M_e ρ^{\delta'}/2; (iv) \|C_s\|_{1,1} ≤ M_e ρ^{\delta'}/2.

Under Assumptions 3.1 and 3.2, the processes \{X_t\} and \{e_t\} satisfying (5) and (6) are zero-mean linear and weakly stationary for any fixed p, and consequently the process \{X_t\} also is. Similar assumptions are made by Forni et al. (2017) in a GDFM context and allow us to control the amount of physical dependence of \{X_t\} across time (Wu and Zaffaroni 2018). Notice that cross-sectional heteroscedasticity of both \{X_t\} and \{e_t\} is allowed for.

Some more comments on Assumption 3.2 are needed. First, the fact that in part (i) there is no dependence on p is natural since B_s has a finite number of columns. Second, part (ii) implies the largest eigenvalue of B_s' is diverging with ρ, because \|B_1\|_2 ≤ \|B_1\|_∞ ≤ M_x ρ^{\delta'}/2, a requirement which is compatible with the idea of pervasive factors at all lags that we impose in Assumption 3.4. Third, part (iii) imposes finite column sums, that is, the \ell∞, norm, for the coefficients C_s, meaning that for any given X_t, the p innovations e_{t₁}, . . . , e_{t_p} have a finite effect for any p ∈ N, which is in agreement with the idiosyncratic nature of \{e_t\} assumed in the GDFM literature. Fourth, to account for some stronger dependence in \{e_t\}, in part (iv) we allow the entire \ell_1 norm of C_s to be diverging with p. This, together with part (iii), implies that the diverging behavior of those coefficients is implicitly due to the row sums, that is, the \ell_1, norm. This is just a useful and natural way of parameterizing the model and, obviously, we could equivalently assume the viceversa or let both row and column sums diverge (compatibly with part (iv)). Fourth, the assumption δ' < α ensures that the common component always dominates the idiosyncratic component when p → ∞.

Now, using the Singular Value Decomposition, we can always write the MA coefficients as \(B_s = U_{L_s}A_{L_s}V_{L_s}, C_s = U_{S_s}A_{S_s}V_{S_s}, s \in Z^+ \cup \{0\}\), where \(U_{L_s}\) is p × r with \(U_{L_s}U_{L_s}^T = I_p\), \(V_{L_s}\) is r × r with \(V_{L_s}V_{L_s}^T = V_{L_s}^TV_{L_s} = I_r\), \(A_{L_s}\) is r × r diagonal, real, and positive definite matrix of singular values. Similarly, \(U_{S_s}\) is p × p with \(U_{S_s}U_{S_s}^T = U_{S_s}^TU_{S_s} = I_p\), \(V_{S_s}\) is p × r with \(V_{S_s}V_{S_s}^T = V_{S_s}^TV_{S_s} = I_r\), \(A_{S_s}\) is p × p diagonal, real, and positive definite matrix of singular values. By means of the following assumption, we impose a low rank plus sparse structure on the filters.

Assumption 3.3. For all s ∈ Z^+ \cup \{0\}: (i) \(U_{L_s} = U_{L_1};\) (ii) \(\|U_{S_s}U_{S_s}^T\|_0 = q_s\), with \(q_s \in N\) and \(q_s < p^2\), for all p ∈ N.

In other words, in part (i) we allow the matrices B'_sB_s to have a different condition number across all lags s, and in part (ii) we allow the matrices C_s'C_s to have a different sparsity pattern across all lags s. Although part (i) might seem restrictive in that it assumes that the space spanned by U_{L_s} is the same at all lags s, it is in fact quite reasonable if, as usually assumed, the spectral density of \{X_t\} has to have rank r at all frequencies (see below). Moreover, we notice that all results in the next section hold locally in a neighborhood of \(L(\theta)\) (see also Chandrasekar, Parrilo, and Wolfsky 2012).

Under Assumption 3.3, the lag-k autocovariances of \{X_t\} and of \{e_t\} are given by \(\Gamma_k(k) = \sum_{n=0}^{\infty} U_{L_s}A_{L_s}A_{L_s+k}U_{L_s}'\) and \(\Gamma_s(k) = \sum_{m=0}^{\infty} U_{S_s}A_{S_s}A_{S_s+k}U_{S_s}'\), respectively, and the spectral density matrices of \{X_t\} and of \{e_t\} are given by \(\Gamma(\theta) = \frac{1}{2\pi} U_{L_s} \left( \sum_{k=-\infty}^{\infty} \sum_{n=0}^{\infty} A_{L_s}A_{L_s+k} e^{-i\theta k} \right) U_{L_s}'\) and \(S(\theta) = \frac{1}{2\pi} \sum_{k=-\infty}^{\infty} \left( \sum_{m=0}^{\infty} U_{S_s}A_{S_s}A_{S_s+k} e^{-i\theta k} \right)\), respectively. Therefore, \(\Gamma_k(k)\) is of rank r for all k ∈ Z and \(L(\theta)\) is of rank r for all \(\theta \in [-\pi, \pi]\). This is natural and standard requirement in GDFM literature (Forni et al. 2000). Moreover, since by Assumption 3.3(ii) \(U_{S_s}\) must be sparse for any \(s \in \mathbb{Z}^+ \cup \{0\}\), then \(\Gamma_s(k)\) will be sparse, and \(S(\theta)\) will also be sparse. In particular, letting \(Q_k = \{i,j = 1, \ldots, p, \Gamma_s(k) \neq 0\}\) with cardinality, say q, and \(Q(\theta) = \{i,j = 1, \ldots, p, \delta(\theta) \neq 0\}\), we have \(Q_k \subseteq \bigcup_{k=0}^{\infty} Q_k\) with cardinality at most q, for some \(\theta \in [-\pi, \pi]\). This has nontrivial implications for the ability of our method to retrieve nonzero entries which are discussed in Remark 5.6.

Summing up, Assumption 3.3 characterizes the spectral density matrix of \{X_t\}, which is \(\Sigma(\theta) = \Lambda(\theta) + S(\theta)\), as having a low rank plus sparse structure at all frequencies \(\theta \in [-\pi, \pi]\).

To fully characterize the low rank property of \(L(\theta)\), we make the following assumption, which strengthens Assumptions 3.2(ii) and 3.3(i).

Assumption 3.4. For all p ∈ N, all \(\theta \in [-\pi, \pi]\), the spectral density \(L(\theta)\) exists, and for all j = 1, . . . , r, there exist continuous functions \(M_{mn}(\theta), M_{mj}(\theta) : [-\pi, \pi] \rightarrow \mathbb{R}^+\) such that \(M_{mn}(\theta) \leq \frac{\lambda_j(L(\theta))}{\rho}\) ≤ \(M_{mj}(\theta),\) with \(\alpha \in (\frac{1}{2}, 1],\) and \(M_{mj}(\theta) \geq M_{mj}(\theta),\) for j = 1, . . . , (r − 1).

The case \(\alpha = 1\) corresponds to the classical assumption in factor model literature of strong pervasiveness of the latent factors (Forni et al. 2017). Here, by allowing the eigenvalues of \(L(\theta)\) to be of order \(\rho^p\), \(\alpha \in (\frac{1}{2}, 1],\) we allow also for weaker factors.

To fully characterize the sparseness property of \(S(\theta)\), we make the following assumption which strengthens Assumptions 3.2(iv) and 3.3(ii).

Assumption 3.5. For all p ∈ N and all \(\theta \in [-\pi, \pi]\), the spectral density \(S(\theta)\) exists and: (i) there exist \(\delta \in [0, \frac{1}{2}]\) and \(\delta_2 > 0,\) such that \(\|S(\theta)\|_0 ≤ \delta_2 p^2\) for all \(\theta \in [-\pi, \pi]\); (ii) there exists a \(M > 0\) such that \(k_p(S(\theta)) > M\).

Part (i) controls the maximum number of nonzeros per row in \(S(\theta)\), which is independent of \(\theta\) because of Assumption 3.3(ii), and it is allowed to increase with p. In particular, the maximum number of nonzeros (across frequencies) in \(S(\theta)\), which we denoted as q, is such that \(q ≤ p\) sup_{\theta \in [-\pi, \pi]} \|S(\theta)\|_0.\) Hence, because of Assumption 3.5, q can grow with p at a maximum rate \(p^{3/2}\). This fully characterizes the sparsity conditions on \(S(\theta)\). Moreover, since Assumption 3.2(iv) implies that
\[ \|S(\theta)\|_1 = O(p^{\delta}), \text{ and since by Assumption 3.5 } \|S(\theta)\|_1 \leq p\|S(\theta)\|_{1,\nu} \leq p\|S(\theta)\|_{0,\nu} \|S(\theta)\|_{\infty} = O(p^{\delta+1}), \text{ then we must have } \delta \leq 1 + \delta. \] In part (ii), by assuming positive definiteness of \(S(\theta)\), we guarantee that also \(\Sigma(\theta)\) is positive definite, because of Weyl's inequality.

**Remark 3.1.** From Assumption 3.5 it follows that the largest eigenvalue of \(S(\theta)\) is at most of order \(p^\delta\), \(\delta \in [0, \frac{1}{2}]\) since \(\|S(\theta)\|_2 \leq \|S(\theta)\|_{\infty} \|S(\theta)\|_{0,\nu} = O(p^\delta)\). Moreover, by Assumption 3.4 the eigenvalues of \(L(\theta)\) are of order \(p^\alpha\), \(\alpha \in \left(\frac{1}{2}, 1\right]\), so that \(\delta < \alpha\). Therefore, there exists an eigenspace between the eigenvalues of \(L(\theta)\) and those of \(S(\theta)\) across all frequencies \(\theta \in [-\pi, \pi]\). In the GDFM setting the presence of an eigenspace growing with \(p\) is the condition which allows for the recovery of the number of factors \(r\) (Hallin and Liška 2007 and Onatski 2009). Notice that here we impose \(\lambda_1(L(\theta)) \approx p^\alpha\), while we just have an upper bound for \(\lambda_1(S(\theta))\), so the eigenspace implied by our assumptions is at least of order \(p^{\alpha-\delta}\), but could be wider.

**Remark 3.2.** It is also worth noticing that in Assumption 3.4 we could easily allow for factors having different degree of pervasiveness, that is, by assuming the eigenvalues at frequency \(\theta\) to be of order \(p^{\bar{\alpha}(\theta)}\) (in that case, \(\alpha\) in Assumption 3.2 would be the infimum over frequencies of all \(\alpha(\theta)\)). Likewise in Assumption 3.5 we could assume \(\|S(\theta)\|_{0,\nu} \leq \bar{\delta}p^{\bar{\delta}(\theta)}\). However, in order to keep the notation simple, we prefer to keep treating \(\alpha\) and \(\delta\) as constants (see also Remark 3.6 for more details).

In order to study the properties of our estimator and guarantee algebraic consistency, we need to formalize further the low rank plus sparsity structure. To this end we introduce the following algebraic matrix varieties for generic integers \(r\) and \(q\):

\[ \mathcal{L}(r) = \{ L \in \mathbb{C}^{p \times p} \mid L \preceq 0, L = UD U^\top, U \in \mathbb{C}^{p \times r}, \] \[ U^\top U = I_r, D \in \mathbb{R}^{r \times r} \text{ diagonal, } \text{rk}(D) = r < p \} \] \[ \mathcal{S}(q) = \{ S \in \mathbb{C}^{p \times p} \mid S \succeq 0, |\text{supp}(S)| \leq q < p^2 \}. \]

where \(\text{supp}(S)\) is the orthogonal complement of \(\text{ker}(S) = \{ v \in \mathbb{C}^p \mid Sv = 0\}\). In other words, \(\mathcal{L}(r)\) is the variety of Hermitian matrices with at most rank \(r\) and \(\mathcal{S}(q)\) is the variety of Hermitian matrices with at most \(q\) nonzero elements. Therefore, under our assumptions \(L(\theta) \in \mathcal{L}(r)\) and \(S(q) \in \mathcal{S}(q)\) for all \(\theta \in [-\pi, \pi]\).

In order to give a rigorous definition of sparsity, we need to introduce further notation. The tangent spaces to \(\mathcal{L}(r)\) and \(\mathcal{S}(q)\) in \(L(\theta)\) and \(S(q)\) are respectively defined, for all \(\theta \in [-\pi, \pi]\), as

\[ T(L(\theta)) = \{ M \in \mathbb{C}^{p \times p} \mid M = U Y_1 + Y_2 U^\top \mid Y_1, Y_2 \in \mathbb{C}^{p \times r}, U \in \mathbb{C}^{p \times p} \text{ diagonal}, \] \[ U^\top U = I_r \} \]

\[ \Omega(S(q)) = \{ N \in \mathbb{C}^{p \times p} \mid |\text{supp}(N)| \subseteq |\text{supp}(S(q))|\}. \]

The following uncertainty principle holds (Chandrasekaran et al. 2011): if \(L(\theta)\) is nearly sparse, \(S(\theta)\) cannot be recovered, and if \(S(\theta)\) is nearly low rank, \(L(\theta)\) cannot be recovered. Therefore, in order to identify \(L(\theta)\) and \(S(\theta)\) at each \(\theta \in [-\pi, \pi]\), we need to ensure that the tangent spaces to \(\mathcal{L}(r)\) and \(\mathcal{S}(q)\) are close to orthogonality. To this end, Chandrasekaran et al. (2011) introduce the following measure for any generic

\[ L \in \mathcal{L}(r) \text{ and } S \in \mathcal{S}(q): \]

\[ \xi(T(L)) = \frac{1}{\max_{M \in T(L)} \|M\|_{\infty}}, \quad \mu(\Omega(S)) = \frac{1}{\max_{N \in \Omega(S)} \|N\|_2}. \]

So if \(\xi(T(L))\) is small then the elements of the tangent space \(T(L)\) are “diffuse,” that is, these elements are not too sparse; as a result \(L\) cannot be very sparse. Similarly, the quantity \(\mu(\Omega(S))\) being small implies that the spectrum of any element of the tangent space \(\Omega(S)\) is “diffuse,” that is, the singular values of \(S\) are not too large. Moreover, (Chandrasekaran et al. 2011, Propositions 3 and 4) show that the following relationships always hold for any generic \(L \in \mathcal{L}(r)\) and \(S \in \mathcal{S}(q)\):

\[ \sqrt{\frac{r}{p}} \leq \xi(T(L)) \leq 2, \]

\[ \min_{1 \leq l \leq p} \sum_{j=1}^p 1(S_{lj} = 0) \leq \mu(\Omega(S)) \leq \max_{1 \leq l \leq p} \sum_{j=1}^p 1(S_{lj} = 0) = \|S\|_{0,\nu}. \]

Moreover, a necessary condition to ensure both parametric and algebraic consistency is

\[ \xi(T(L(\theta))) \cdot \mu(\Omega(S(\theta))) \leq \frac{1}{54}, \quad \theta \in [-\pi, \pi], \]

which guarantees that \(L(\theta)\) is far from sparsity and \(S(\theta)\) is far from rank-deficiency. Indeed, the smaller is the product between the dual norms \(\xi(T(L(\theta)))\) and \(\mu(\Omega(S(\theta)))\), the closer the two spaces \(\mathcal{L}(r)\) and \(\mathcal{S}(q)\) are to orthogonality, thus, making easier to perform low rank plus sparse decomposition effectively.

To control the rank-sparsity incoherence measures we make the following assumption, which guarantees that (9) and (10) are satisfied.

**Assumption 3.6.** For all \(p \in \mathbb{N}\), there exist \(\kappa_L, \kappa_S > 0\) with \(\sqrt{\kappa_L} \leq 1, \sqrt{\kappa_S} \leq 1\), \(\xi(T) \leq \frac{\sqrt{\kappa_L}}{54}\) and \(\mu(\Omega) \leq \kappa_S p^\delta\), for all \(\theta \in [-\pi, \pi]\), and where \(\delta\) and \(\delta\) are defined in Assumption 3.5.

Now, because of Assumption 3.3, \(T(L(\theta))\) depends only on the rank of \(L^s = \Gamma_s(0)\) which is \(r\) and, \(\Omega(S(\theta))\) depends only on the support of \(S_{\infty} = \sum_{\epsilon = -\infty}^{\infty} \Gamma_{\epsilon}(k)\), which has dimension \(q\). For this reason, hereafter, we use the shorthand notation \(T, \Omega, \xi(T), \mu(\Omega)\).

**Remark 3.3.** The validity of our setup goes beyond the VMA(\(\infty\)) framework discussed so far. First, it is obvious that any stable VARMA with finite lags would fit into our framework. We refer to Appendix B, supplementary materials for specific conditions on the VARMA coefficients such that our assumptions are satisfied. Second, as long as \(\{u_t\}\) and \(\{e_t\}\) can be expressed as measurable functions of iid processes satisfying Assumption 3.1, it can be shown that the theory developed in this article is still valid. We refer to Appendix C, supplementary materials for details and the extension to the nonlinear case in which \(\{u_t\}\) and \(\{e_t\}\) are allowed to be conditionally heteroscedastic.
4. Estimation

Suppose now to observe a sample of $p$-dimensional data vectors with size $T$, that is, we observe $\{X_{i,t}, i = 1, \ldots, p, t = 1, \ldots, T\}$, and we compute the ALSF $\hat{\Sigma}(\theta) = \hat{L}(\theta) + \hat{S}(\theta)$ such that it satisfies (1) which, for convenience, we rewrite here:

$$\hat{L}(\theta), \hat{S}(\theta) = \arg \min_{\hat{L}(\theta) \geq 0, \hat{S}(\theta) \geq 0} \{ \frac{1}{2} \| \hat{\Sigma}(\theta) - (\hat{L}(\theta) + \hat{S}(\theta)) \|_F^2 + \psi \| \hat{L}(\theta) \|_s + \rho \| \hat{S}(\theta) \|_1 \} , \quad (11)$$

where $\hat{\Sigma}(\theta)$ is the smoothed periodogram estimator defined in (3), $\hat{L}(\theta)$ and $\hat{S}(\theta)$ indicate generic values of the matrices, and $\psi$ and $\rho$ are positive threshold parameters and their choice is discussed in Section 7.

The minimization problem (11) is a non-smooth convex optimization problem which is the tightest convex relaxation of the following NP-hard problem:

$$\min_{\hat{L}(\theta), \hat{S}(\theta) \geq 0} \frac{1}{2} \| \hat{\Sigma}(\theta) - (\hat{L}(\theta) + \hat{S}(\theta)) \|_F^2 + \psi \| \hat{L}(\theta) \|_s + \rho \| \hat{S}(\theta) \|_1$$

which would be the natural target under the low rank plus sparse assumption. Indeed, we know that: (i) $\| \hat{\Sigma}(\theta) \|_1$ is the tightest convex relaxation of $\{ \| \hat{\Sigma}(\theta) \|_1 \}$ (Donoho 2006); (ii) $\| \hat{L}(\theta) \|_s$ is the tightest convex relaxation of $\text{rank}(\hat{L}(\theta))$ (Fazel, Hindi, and Boyd 2001).

In practice, the solution of (11) is computed as follows. For any given frequency $\theta_h = \frac{2\pi h}{MT}$, with $|h| \leq M_T$, we apply the following iterative procedure:

1. set $(L_0(\theta_h), S_0(\theta_h)) = \left( \text{diag}(\hat{\Sigma}(\theta_h)), \text{diag}(\hat{\Sigma}(\theta_h)) \right)$, $\eta_0 = 1$, and initialize $Y_0(\theta_h) = L_0(\theta_h)$ and $Z_0(\theta_h) = S_0(\theta_h)$;

2. for $k \geq 1$, repeat:
   a. compute
      $$\frac{1}{2} \| Y_{k-1}(\theta_h) + Z_{k-1}(\theta_h) - \hat{\Sigma}(\theta_h) \|^2_F = \frac{1}{2} \| Y_{k-1}(\theta_h) + Z_{k-1}(\theta_h) - \hat{\Sigma}(\theta_h) \|^2_F$$
      $$\text{sgn} ( S_{k-1}(\theta_h) ) = Y_{k-1}(\theta_h) + Z_{k-1}(\theta_h) - \hat{\Sigma}(\theta_h)$$
   b. apply the singular value thresholding operator of Cai, Candès, and Shen (2010) $T_{\psi}(\cdot)$ to $E_{Y_k}(\theta_h) = Y_{k-1}(\theta_h) - \frac{1}{2} ( Y_{k-1}(\theta_h) + Z_{k-1}(\theta_h) - \hat{\Sigma}(\theta_h) )$ and set $L_k(\theta_h) = T_{\psi}(E_{Y_k}(\theta_h));$
   c. apply the soft-thresholding operator of Daubechies, Defrise, and De Mol (2004) $T_{\psi}(\cdot)$ to $E_{Z_k}(\theta_h) = Z_{k-1}(\theta_h) - \frac{1}{2} ( Y_{k-1}(\theta_h) + Z_{k-1}(\theta_h) - \hat{\Sigma}(\theta_h) )$ and set $S_k(\theta_h) = T_{\psi}(E_{Z_k}(\theta_h));$
   d. set $(Y_k(\theta_h), Z_k(\theta_h)) = (L_k(\theta_h), S_k(\theta_h)) + \sum_{k=0}^{n_k} \left[ (L_{k}(\theta_h), S_{k}(\theta_h)) - (L_{k-1}(\theta_h), S_{k-1}(\theta_h)) \right]$ where $n_k = \left\lceil \frac{1}{4} \log_2 \frac{1}{\epsilon} \right\rceil$;
   e. stop if $\left\lceil \frac{1}{4} \log_2 \frac{1}{\epsilon} \right\rceil + \frac{1}{\epsilon} \leq \epsilon$, where $\epsilon$ is a prescribed precision level (we set $\epsilon = 0.01$);  
3. set $\hat{L}(\theta_h) = Y_k(\theta_h)$ and $\hat{S}(\theta_h) = Z_k(\theta_h)$.

The two thresholding operators introduced in the above algorithm are defined as follows.

(i) Singular value thresholding operator: let the Singular Value Decomposition of a positive semidefinite complex symmetric matrix $A$ be $A = U_A \Lambda_A U_A^T$, then, define $T_{\psi}(A) = U_A \Lambda_{\psi,A} U_A^T$, where $\Lambda_{\psi,A}$ is a diagonal matrix with $\psi$th diagonal element $\Lambda_{\psi,A} = \max (0, \Lambda_A - \psi)$.

(ii) Soft-thresholding operator: for a positive definite complex symmetric matrix $A$ define $T_{\psi}(A) = \frac{\Lambda_A}{(\Lambda_A + \psi)^{1/2}} \max (\Lambda_A, \Lambda_A - \psi)^{1/2} - \rho, 0)$.

Two features of ALSE must be stressed. First, not only ALSE produces the estimates $\hat{L}(\theta)$ and $\hat{S}(\theta)$, but it also produces estimates of $\text{rank}(\hat{L}(\theta)) = r$ and of $\text{supp}(\hat{S}(\theta))$ and therefore of its cardinality $q$. Second, the solution of the above minimization can be searched without the need of constraining $\hat{L}(\theta)$ and $\hat{S}(\theta)$ to the manifolds $L(r)$ and $S(q)$ defined in (7–8) (see also Remark 5.6 and Appendix A, supplementary materials).

5. Consistency

We now prove the algebraic and parametric consistency of the pair of estimates $(\hat{L}(\theta), \hat{S}(\theta))$, and, in order to do this, we introduce two definitions, taken from Chandrasekaran, Parrilo, and Willsky (2012). First, we say that $(\hat{S}(\theta), \hat{L}(\theta))$ is algebraically consistent if the following conditions hold with probability tending to one, for any given $\theta \in [-\pi, \pi]$; (i) $\{ \hat{L}(\theta) \} = \text{rank}(\hat{L}(\theta))$; (ii) $\text{sgn}(\text{Re}(\hat{S}(\theta)))$ = $\text{sgn}(\text{Re}(S_0(\theta)))$, for all $i, j = 1, \ldots, p$ (by convention we let $\text{sgn}(0) = 0$); (iii) $\{ \hat{L}(\theta) \}$ and $\{ \hat{S}(\theta) \}$ are positive definite and $\hat{L}(\theta)$ is positive semidefinite. Condition 2 is also often referred to as sparsity consistency (Lam and Fan 2009). Second, we use the classical definition of parametric consistency, which holds if the estimates $(\hat{S}(\theta), \hat{L}(\theta))$ are close to $(S(\theta), L(\theta))$, for any given $\theta \in [-\pi, \pi]$, with high probability in $\ell_2$ norm.

In order to prove consistency of the pre-estimator of the spectral density matrix $\hat{\Sigma}(\theta)$, defined in (3), we make the following standard assumption on the kernel function and its bandwidth.

Assumption 5.1. The kernel function $K(\cdot)$ is even, bounded, with support $[-1, 1]$, and bandwidth $M_T$, such that: (i) for some $k > 0$, $|K(k) - 1| = O(s^k)$, as $s \to 0$; (ii) $\int_{-\infty}^{\infty} K^4(s)ds < K$ for some finite $K$; (iii) $\sum_{r \in \mathbb{Z}} \lambda_{s} < 0 < \zeta$ and $s \to 0$; (iv) $c_1 \zeta s \leq M_T \leq c_2 s^k$, for some $c_1, c_2 > 0$ and $0 < s < \zeta < -1 < 1 < 2(k + 1)$.

By properly adapting the results of Wu and Zaffaroni (2018) and Zhang and Wu (2021) to the intermediate spikiness-sparsity regimes described in Section 3, we prove uniform consistency over frequencies of the smoothed periodogram pre-estimator (3).

Lemma 5.1. Under Assumptions 3.1, 3.2, and 5.1, there exists $G, G' > 0$ and $d > 4$, independent of $p$ and $T$, such that, as $T \to \infty$ and for all $p \in \mathbb{N}$, for $\theta_h = \frac{2\pi h}{MT}$, $|h| \leq M_T$:

1. $Pr \left( \max_{|h| \leq M_T} \frac{1}{p} \| \hat{\Sigma}(\theta_h) - \Sigma(\theta_h) \|_2 \leq G \sqrt{\frac{M_T \log(M_T)}{p}} \right) \to 1$;
2. $Pr \left( \max_{|h| \leq M_T} \| \hat{\Sigma}(\theta_h) - \Sigma(\theta_h) \|_\infty \leq G' \max \left( \frac{M_T \log(M_T)}{p}, \sqrt{\frac{M_T \log(M_T)}{p}} \right) \right) \to 1$. 

The proof of part 1 is new to this article and it is a nontrivial generalization to the case of weakly varying factors of the result proved in Forni et al. (2017, Proposition 6) derived for the GDFM when \( \alpha = 1 \). Part 2 is derived by adapting the results of Zhang and Wu (2021, Proposition 4.3) to the present context.

Remark 5.1. Notice that in part 1 of Lemma 5.1 the bias term, which is of order \( \frac{1}{M_T^2} \), is not included, since, for all \( M_T \) satisfying Assumption 5.1(iv), this term is always dominated by the variance term. Indeed, while the optimal choice of expanding variance and squared bias is \( M_T = O(T^{1/(2k+1)}) \), we are instead assuming \( M_T = O(T^\xi) \) with \( \xi > \frac{1}{k+1} \), and with this choice of \( M_T \) the mean squared error of the smoothed periodogram is dominated by the variance, while the squared bias becomes negligible, as \( T \to \infty \). Typical values of \( k \) are 1 if we choose the Bartlett kernel, or 2 if we choose the Parzen kernel.

The following theorem, proving consistency of ALSE, is the main result of the article.

Theorem 5.1. Define \( \varphi_{p,T,d} = \max(M^2 p^{-d} \log p \varepsilon^2, \sqrt{M^2 \log(M^2 p^d)}) \), with \( d = \max(d_a, d_\epsilon) > 4 \), and define \( \psi_0 = \frac{\psi}{\kappa^{\mu}} \) and \( m_p = \sup_{\theta \in [-\pi, \pi]} |S(\theta)|_0, \) Set \( \psi_0 = \varphi_{p,T,d} \psi \), and \( \psi, \kappa \in [\varphi(T), \frac{1}{\varphi(T)}] \). Suppose that Assumptions 3.1–3.6 and 5.1 hold and suppose also that for all \( \theta \in [-\pi, \pi] \) there exist \( G_2, G_3 > 0 \) such that: (I) \( \lambda_1(L(\theta)) > G_2 \psi \), (II) \( |S(\theta)|_{\text{lim. min.}} \leq G_3 \psi \). Then, there exist positive reals \( \kappa_1 \) and \( \kappa_2 \) independent of \( p \) and \( T \) such that, as \( T \to \infty \) and for all \( p \in \mathbb{N} \), the following hold:

1. \( \mathcal{P} \left( \max_{|\theta| \leq M_T} \frac{1}{\kappa^{\mu}} \| \hat{S}(\theta) - L(\theta) \|_2 \leq \kappa_1 \psi_{p,T,d} \right) \to 1; \)
2. \( \mathcal{P} \left( \max_{|\theta| \leq M_T} \| \hat{S}(\theta) - S(\theta) \|_\infty \leq \kappa_2 \psi_{p,T,d} \right) \to 1; \)
3. \( \mathcal{P} \left( \frac{\text{rk}(L(\theta))}{T} \to r \right) \to 1, \) for all \( \theta \in \mathbb{R} \), \( |\theta| \leq M_T; \)
4. \( \mathcal{P} \left( \text{sgn}(\text{Re}(\hat{S}(\theta)_{ij})) - \text{sgn}(\text{Re}(S(\theta)_{ij})) \right) \to 1, \) for all \( \theta \in \mathbb{R}, \) \( |\theta| \leq M_T \) and all \( i, j = 1, \ldots, p; \)
5. \( \mathcal{P} \left( \max_{|\theta| \leq M_T} \| \hat{S}(\theta) - S(\theta) \|_2 \leq \kappa_2 m_p \psi_{p,T,d} \right) \to 1; \)
6. \( \mathcal{P} \left( \max_{|\theta| \leq M_T} \frac{1}{\kappa^{\mu}} \| \hat{S}(\theta) - S(\theta) \|_2 \leq \kappa_1 \psi_{p,T,d} \right) \to 1. \)

This theorem is a generalization of the results by Chandrasekaran, Parrilo, and Willsky (2012) for the covariance matrix of independent data having a low-rank plus sparse structure. In particular, here we go one step further by considering the case of spectral densities with \( r \) ultra-spiking eigenvalues in the sense that they increase with \( p^{\alpha}, \alpha \in (\frac{1}{2}, 1), \) as prescribed by Assumption 3.4, and where the residual component has a number of nonzeros also growing with \( p^{\delta}, \delta \in [0, \frac{1}{2}], \) as prescribed by Assumption 3.5.

The following remarks provide more intuitions about our results and a comparison with the existing literature.

Remark 5.2. While under the assumptions of Chandrasekaran, Parrilo, and Willsky (2012) the eigen-gap does not depend on \( p \), ours is widening as \( p \) increases. This latter assumption, which is standard in GDFM literature (Forni et al. 2017) and other factor model works (Fan, Liao, and Mincheva 2013), makes identification and thus disentangling of the low rank and sparse component easier, but on the other hand it implies convergence rates that depend polynomially on \( p \), which appear in the rescaling terms in parts 1 and 6, and are discussed in the next remarks.

Remark 5.3. Because of the spiking behavior of the \( r \) eigenvalues of \( L(\theta) \), it is natural to work with the minimization (11) rescaled by \( p^{\alpha} \). This implies that we obtain convergence rates for \( \hat{L}(\theta) \) rescaled by \( p^{\alpha} \). This is standard in this kind of literature, see for example, Fan, Liao, and Mincheva (2013, Theorem 3) for the case of covariance matrices. This explains also the reason why we give conditions for \( \psi_0 = \frac{\psi}{\kappa^{\mu}} \) rather than for \( \psi \).

Remark 5.4. The choice of \( \psi_0 \) to be increasing with \( p^{\delta} \) and of \( \gamma \) to be decreasing with \( p^{\delta} \), because of the definitions of \( \xi(T) \) and \( \mu(\Omega) \) in Assumption 3.6, implies that \( \rho_0 = \gamma \psi_0 \) does not depend on \( \delta \). This choice makes sense, indeed the sparsity of \( S(\theta) \) is at most \( O(p^{\beta}) \) and, therefore, the less sparse is \( S(\theta) \) (higher \( \delta \)) the more this component will become important and the more weight we need to attach to the penalty \( \| L(\theta) \|_2 \) in (11). Such penalization scheme gives rise to the term \( p^{\delta} \) in parts 1 and 6. The rates of convergence are slower with respect to those in Lemma 5.1 obtained for the pre-estimator \( \hat{S}(\theta) \). This is the price to pay for considering a penalized estimator.

Related to this, it is also useful to notice that parts 1 and 2 could be restated as follows: \( \mathcal{P} \left( \max_{|\theta| \leq M_T} \frac{1}{p^{\alpha}} \| \hat{L}(\theta) - L(\theta) \|_2 \leq \kappa \psi_0 \right) \to 1 \) and \( \mathcal{P} \left( \max_{|\theta| \leq M_T} \| \hat{S}(\theta) - S(\theta) \|_\infty \leq \kappa \rho_0 \right) \to 1, \)

Remark 5.5. Part 5 shows that the estimation error of \( \hat{S}(\theta) \) depends on the sparsity of \( S(\theta) \) as measured by \( m_p \). The smaller \( m_p \) is (more sparse), the smaller the estimation error is. By noticing that, because of Assumptions 3.3 and 3.5, it holds that \( m_p \leq \delta_2 p^{\delta} \), this result could also be restated as \( \mathcal{P} \left( \max_{|\theta| \leq M_T} \frac{1}{p^{\alpha}} \| \hat{S}(\theta) - S(\theta) \|_2 \leq \kappa_2 m_p \psi_{p,T,d} \right) \to 1, \) which represents the worst-case scenario in which \( S(\theta) \) is the least sparse compatibly with Assumption 3.5. Notice also that, since \( \delta < \alpha \), the effect of sparsity in part 6 is always dominated by the estimation error of the low rank component \( \hat{L}(\theta) \).

Remark 5.6. For the estimation algorithm to work properly it is crucial that \( \text{rk}(L(\theta)) \) is constant across frequencies as required by Assumption 3.3(i). As already mentioned this is standard in GDFM literature. Moreover, we notice that, although in principle we allow for the number of nonzero elements of \( S(\theta) \) to be frequency dependent, ALSE is in fact able to control only for the maximum number of nonzeros across all frequencies. This is because, as shown in Appendix A, supplementary materials, ALSE is equivalent to looking for a solution of the minimization (11) restricted to \( \Omega \), which depends only on \( S_{\infty} \), and in turn it holds that \( \sup_{\theta \in [-\pi, \pi]} |S(\theta)|_0 \leq \| S_{\infty} \|_{0,v} \).

By looking at the rate \( \varphi_{p,T,d} \), it is clear that in general \( p \) can grow at most polynomially in \( T \), so let us assume \( p \approx T^q \) for
some \( \eta > 0 \). Let us also consider the standard case in which we use the Bartlett kernel with bandwidth to \( M_T = \sqrt{T} \) to compute the pre-estimator, that is, we set \( \zeta = \frac{1}{3} \) and \( k = 1 \) in Assumption 5.1. These are the same choices used in Section 8 and in the simulation study. Furthermore, recall that \( d \) denotes the minimum number of moments of \( \{ u_t \} \) and \( \{ e_t \} \) we require to exist, which by Assumption 3.1 must be such that \( d > 4 \).

Clearly, under these conditions, the second term in \( \varphi_{p,T,d} \) is always decreasing as \( T \to \infty \). A necessary condition for the first term to converge is \( d > 2(1 + \eta) \), which is nonbinding as long as \( \eta \leq 1 \) since, by assumption, \( d > 4 \). However, if \( p \) increases faster than \( T \), that is, \( \eta > 1 \), we must guarantee the existence of more moments to still have consistency. If, for example, \( \eta = 2 \), then we must have \( d > 6 \). Summing up, in general, the larger is \( p \) the lighter the tails have to be in order to guarantee consistency.

As far as the rate of consistency is concerned, we notice that if \( d > 8(1 + \eta) \) we get the classical consistency rate \( \varphi_{p,T,d} = \sqrt{\frac{M_T \log(M_T)}{p}} \), which is similar to results for the estimation of large covariance matrices, see also Remark 5.7. In particular, this rate is achieved, regardless of \( \eta \), in the sub-Gaussian case, that is, when we can set \( d = \infty \), in which case we can also allow for larger values of \( p \) as long as \( \log p = o(T) \). However, in general, when we are in presence of heavy tails, we would typically have the slower consistency rate \( \varphi_{p,T,d} = \frac{M_T \sqrt{\log\log p}}{p^{1 + \delta}} \), unless \( p \) is very small.

To have consistency we also need conditions (I) and (II) to hold. Condition (I) must be compatible with Assumption 3.4, which requires \( \lambda_\alpha(\Sigma(\theta)) < p^\alpha \) for all \( \alpha \in [-\pi, \pi] \). Thus, we must require that

\[
\frac{\eta(\theta_0)}{\zeta \log T} \to 0, \quad \text{as } p, T \to \infty.
\]

In light of the previous comments, setting again \( M_T = \sqrt{T} \), we must have either \( \frac{\sqrt{\log\log p}}{\sqrt{T} \log \log T} \to 0 \) or \( \frac{\log \log p}{\log T} \to 0 \), in the heavy tail or in the sub-Gaussian case, respectively. Both conditions are always true when \( \theta_0 = 0 \), while they require larger \( T \) the smaller the eigengap gets (decreasing \( \alpha - \delta \)). A sufficient condition for both to hold is \( \delta \leq \frac{\eta - 1/3}{\eta} + \epsilon \), for some \( \epsilon > 0 \).

Condition (II) must be compatible with the obvious requirement \( \| S(\theta) \|_{\min, \ell^1, \ell^2} \leq G_4' \) for some \( G_4' > 0 \). This is always true, indeed, by Assumption 3.6, and condition (II) can be written as

\[
\| S(\theta) \|_{\min, \ell^1, \ell^2} \geq G_3 \frac{k_\psi \varphi_{p,T,d}}{\sqrt{\kappa_0}},
\]

which under the conditions on \( p, T, \) and \( d, \) stated above, is always decreasing to zero as \( p, T \to \infty \).

The following Corollary characterizes the inverse of the estimated spectral density matrix.

**Corollary 5.1.** Suppose that all assumptions of Theorem 5.1 hold, and suppose also that \( m_p \varphi_{p,T,d} \to 0 \), as \( T \to \infty \) and for all \( p \in \mathbb{N} \). Then, there exist positive reals \( k_5 \) and \( k_4 \) independent of \( p \) and \( T \) such that, as \( T \to \infty \) and for all \( p \in \mathbb{N} \), the following hold:

1. \( \mathcal{P} \left( \min_{| \theta | \leq M_T} \lambda_\alpha(\hat{\Sigma}(\theta)) \geq k_5 \right) \to 1; \)
2. \( \mathcal{P} \left( \max_{| \theta | \leq M_T} \| \hat{\Sigma}^{-1}(\theta) - S^{-1}(\theta) \|_2 \leq k_4 m_p \varphi_{p,T,d} \right) \to 1; \)
3. \( \mathcal{P} \left( \min_{| \theta | \leq M_T} \lambda_\alpha(\hat{\Sigma}(\theta)) \geq k_3 \right) \to 1. \)

If we further suppose that \( p^{\alpha + 3} m_p \varphi_{p,T,d} \to 0 \), then, there exists a positive real \( k_5 \) independent of \( p \) and \( T \) such that, as \( T \to \infty \) and for all \( p \in \mathbb{N} \), the following hold:

4. \( \mathcal{P} \left( \max_{| \theta | \leq M_T} \frac{1}{p m_p} \| \hat{\Sigma}^{-1}(\theta) - S^{-1}(\theta) \|_2 \leq k_5 \varphi_{p,T,d} \right) \to 1. \)

**Remark 5.7.** The results in Theorem 5.1 and Corollary 5.1 can be compared to those obtained by Fan, Liao, and Mincheva (2013, Remark 2) for the estimation of a large covariance matrix with a low rank component generated by \( r \) strong factors, plus a sparse component having bounded eigenvalues and \( \ell_1 \) norm equal to \( m_p \) (by considering the exact sparse case in their work, i.e., when setting \( q = 0 \) therein). They also assume to observe time series drawn from a distribution with sub-exponential tails, thus, assuming all moments to exist, and obtain for their estimator of the sparse component and its inverse a consistency rate \( O_p \left( m_p \sqrt{\log p} \right) \). This is also the same rate obtained by Bickel and Levina (2008, Theorem 2) when considering the case of a large purely sparse covariance matrix with Gaussian entries. Using our notation the assumptions of Fan, Liao, and Mincheva (2013) correspond to setting \( \alpha = 1, \delta = 0, \) and \( d = \infty \), in which case from part 5 of Theorem 5.1 and from part 2 of Corollary 5.1, we would get a similar rate, which, with our notation, is \( O_p \left( m_p \sqrt{\frac{M_T \log(M_T)}{p}} \right) \). Concerning the estimator of the entire covariance matrix, Fan, Liao, and Mincheva (2013, Remark 3) obtain a consistency rate \( O_p \left( \frac{m_p \sqrt{\log(M_T/p)}}{p} \right) \), which is similar to the rate \( O_p \left( \frac{m_p \sqrt{M_T \log(M_T/p)}}{p} \right) \) that we would get from part 6 of Theorem 5.1 when \( \alpha = 1, \delta = 0, \) and \( d = \infty \).

**Remark 5.8.** Concerning the estimator of the inverse spectral density in part 4 of Corollary 5.1, we obtain the same rate of the estimator of \( \Sigma(\theta) \) given in part 6 of Theorem 5.1, which is worse than what obtained for the inverse covariance matrix by Fan, Liao, and Mincheva (2013). To obtain a comparable rate we should follow the approach of Fan, Liao, and Mincheva (2013, eq. (2.13) and Appendix C.4.2) and define the inverse estimator by explicitly exploiting the existence of a low-rank component and then applying the Sherman-Morrison-Woodbury formula. This, however, would require to explicitly estimate also the coefficients \( B_5 \) of the common filters (see (5)), which is beyond the scope of this article, and, therefore, it is left for further research.

6. Unshrinking

ALSE may suffer from systematic sub-optimality for what concerns the estimated eigenvalues. Indeed, if \( p \) is large and the latent eigenvalues are spiked (\( \alpha \approx 1 \)), the singular value thresholding procedure by Cai, Candès, and Shen (2010) may lead to the over-shrinkage of the eigenvalues of \( L(\theta_0) \). For this reason, following Farné and Montanari (2020), we propose to un-shrink those eigenvalues. The new estimate of \( S(\theta) \) is then obtained by keeping fixed the off-diagonal sparsity pattern recovered, and deriving its diagonal by difference from the diagonal of \( \Sigma(\theta_0) \). More specifically, let \( \hat{\tau} = \kappa_0 \hat{L}(\theta_0) \) and consider the spectral decomposition \( \hat{L}(\theta_0) = W(\theta_0) D(\theta_0) W^T(\theta_0) \), where \( D(\theta_0) \) is \( \hat{\tau} \times \hat{\tau} \) diagonal matrix of the nonzero eigenvalues of \( L(\theta_0) \), and \( W(\theta_0) \) is the \( p \times \hat{\tau} \) matrix of corresponding normalized
In solving problem (11), the choice of the thresholds \( \psi \) and \( \rho \) is a nontrivial issue. Let us denote the solutions of (11) with \( L \) rows of the \( p \times p \) matrix \( \Sigma(\theta_h) \), where

\[
\psi(\theta_h) = \frac{\Sigma(\theta_h)}{\sqrt{\text{tr}(\Sigma(\theta_h))}}, \quad \rho(\theta_h) = \frac{\text{tr}(\Sigma(\theta_h))}{\text{tr}(\Sigma(\theta_h))},
\]

then the algebraic consistency of ALSE, proved in parts 3 and 4 of Theorem 5.1, is preserved by UNALSE. Theorem A.1 in Appendix A also shows that, under the same assumptions and conditions of Theorem 5.1, UNALSE is the closest (according to the Frobenius norm) algebraically consistent estimator to the smoothed periodogram pre-estimator.

7. Threshold Selection

In solving problem (11), the choice of the thresholds \( \psi \) and \( \rho \) is a nontrivial issue. Let us denote the solutions of (11) with \( L \) rows of the \( p \times p \) matrix \( \Sigma(\theta_h) \), where \( L \) is the top eigenvalue shows a maximum at \( \hat{\theta} \), e.g., McCracken and Ng (2020). Throughout, we compute the smoothed periodogram by setting \( M_T = \sqrt{T} \), defined as in (13) and using the Bartlett kernel. UNALSE is computed as described in Sections 4 and 6, while the thresholds are chosen as explained in Section 7. Results are reported using frequencies \( f_h = \frac{\theta_h}{2\pi M_T}, |\theta_h| \leq M_T \), measured in cycles per unit of time. Thus, for a given frequency \( f_h \) the corresponding period is given by \( f_h^{-1} \) and, since our data is quarterly, its unit of measure are quarters, with one quarter being equal to three months.

In the first panel of Figure 1 we show the four largest eigenvalues, rescaled by \( \hat{\rho} \), of the smoothed periodogram estimator. The top eigenvalue shows a maximum at \( \hat{\theta} \) and a second peak at \( \hat{\theta} = 0.31 \), corresponding to periods of about 3.5 years and 9 months, respectively. Note that 3.5 years is around the typical period of a business cycle, while the higher frequency peak is typically more related to nominal variables such as inflation (see also the results in Barigozzi and Luciani, 2021, and below). The estimated rank by UNALSE is \( \hat{r} = 2 \) at all frequencies, which is an agreement with the recent findings by Avarucci et al. (2022), who interpret the two common factors as a demand (high-frequency) and a supply (low-frequency) factors (see also Angelots, Collard, and Dellas 2020). The second panel of Figure 1 shows the proportion of latent variance \( \rho(f_h) = \frac{\text{tr}(\hat{L}(f_h))}{\text{tr}(\Sigma(f_h))} \) so the contribution of \( \hat{L}(f_h) \), which follows the pattern of the leading eigenvalues of \( \hat{\Sigma}(f_h) \), hence, it captures the business cycle frequency. The third panel of Figure 1 reports the proportion of residual covariance \( \hat{\Sigma}(f_h) = \sum_{i=1}^{p} \sum_{j=1}^{p} |\hat{\beta}_{ij}(f_h)| \).
summarizing the contribution of \( \hat{S}(f_h) \), which is maximum at frequency \( f_h = 0.48 \), corresponding to a period of about 6 months. Secondary maxima are at \( f_h = 0.03 \), that is, a period of 8 years, and \( f_h = 0.24 \), corresponding to a period of 1 year. Finally, the fraction of nonzeros in the fourth panel of Figure 1 has a similar pattern.

In Figure 2 (first two panels) we show heat-maps of \( \hat{L}(f_h) \) at frequencies 0.07 and 0.31. The variables contributing more to the comovements at the low frequency \( f_h = 0.07 \) are either related to the labor market (variables from \( i = 70, \ldots, 94 \)) such as (i) Unemployment rate; (ii) Average Mean Duration of Unemployment; (iii) All Employees in various sectors, or are related to credit market and the real economy (variables from \( i = 25, \ldots, 37 \)) such as GDP, Investment, and Consumption growth rates. All these are the variables typically driving the business cycle. Note that these variables drive also the comovements at \( f_h = 0 \). At higher frequency \( f_h = 0.31 \) the most relevant variables are the growth rates, that is, inflation, of Consumer and Production Price Indexes (variables from \( i = 7, \ldots, 19 \)). Thus, the two main sources of common variation are related to: (i) the real economy, and (ii) to the nominal economy.

In Figure 2 (last two panels) we show heat-maps of \( \hat{S}(f_h) \) at frequencies 0.03 and 0.48. It is worth mentioning some prominent co-spectral relationship at selected frequencies. At frequency \( f_h = 0.03 \), the following pairs display strong co-dependence: (i) Consumer Loans at All Commercial Banks and Total Consumer Credit (Owned and Securitized); (ii) 3-Year Treasury Constant Maturity Rate and 10-Year Treasury Constant Maturity Rate; (iii) Compensation Per Hour in the Manufacturing Sector and in the Business Sector. At frequency \( f_h = 0.24 \), the strongest co-spectral relationships are instead (not shown): (i) Producer Price Index of Industrial Commodities and of All Commodities; (ii) Producer Price Index of Industrial Commodities and Treasury Constant Maturity rate. These local, co-movements can be interpreted as due to weaker, idiosyncratic, factors related to the banking system represented by credit variables (low-frequency), to prices and interest rates, that is, related to monetary policy (mid-frequency), and to the trade/consumption dimension of the U.S. economy (high-frequency). No relevant idiosyncratic comovements are found in the long run, that is, at \( f_h = 0 \).

9. Conclusions

In this article, we consistently estimate the spectral density matrix under the assumption of a low rank plus sparse multivariate spectrum for the data. We prove that the nuclear norm plus \( l_1 \) norm heuristics consistently recovers across frequencies the spectral components and their sum, as well as the dynamic rank and the residual sparsity pattern. We call the resulting estimators UNALSE (UNshrunk ALgebraic Spectral Estimator). The empirical implications of our approach are discussed on a standard U.S. macroeconomic dataset, showing that we are able to catch the driving variables of the latent dynamics, as well as the particular strength of specific local relationships which might be heterogeneous across frequencies. This opens up the way to enhanced dynamic factor scores estimation and temporal network analysis.
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The authors report there are no competing interests to declare.

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