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Estimation of Common Factors for Microstructure Noise and Efficient Price in a High-frequency Dual Factor Model

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Abstract

We develop the Double Principal Component Analysis (DPCA) based on a dual factor structure for high-frequency intraday returns data contaminated with microstructure noise. The dual factor structure allows a factor structure for the microstructure noise in addition to the factor structure for efficient log-prices. We construct estimators of factors for both efficient log-prices and microstructure noise as well as their common components, and provide uniform consistency of these estimators when the number of assets and the sampling frequency go to infinity. In a Monte Carlo exercise, we compare our DPCA method to a PCA-VECM method. Finally, an empirical analysis of intraday returns of S&P 500 Index constituents provides evidence of co-movement of the microstructure noise that distinguishes from latent systematic risk factors.

Keywords: Cointegration, Factor model, High-frequency data, Microstructure noise, Non-stationarity, *AMS Subject Classification*:

1. Introduction

Factor models are widely used in many scientific fields, and in particular in the study of financial data. Their popularity is partly due to the easiness of their implementation and their effectiveness in dimension reduction. More and more observable factors have been investigated and reported (see. e.g., Ross (1976), Sharpe (1994), Fama and French (1993, 2015) and Carhart (1997)) as driving stock returns. Researchers have also found common components in other attributes of financial assets such as volatility and liquidity. For example, Chordia et al. (2000) document the commonality in liquidity, which remains significant after controlling for volatility, volume, and price. The factor structure is not found in isolation. Indeed, price and other attributes of stocks have been found to have correlated common factors. Hasbrouck and Seppi (2001) use principal component analysis to show that common factors exist in order flows and equity returns. In addition, using canonical correlation analysis, they find that the common factor in returns is highly correlated with the common factor in order flows. Hallin and Liška (2011) propose a two-step general dynamic factor method to account for a joint factor structure of sub-panels, which is further developed by Barigozzi and Hallin (2016) and Barigozzi and Hallin (2017) for extracting the market volatility shocks. They find that returns and volatilities can be decomposed into four mutually orthogonal components: a strongly idiosyncratic component, a strongly common component, a weakly common component, and a weakly idiosyncratic component.

The increasing availability of high-frequency transaction data motivates applying this methodology to intraday stock prices. However, this setting raises certain theoretical and computational challenges. Compared to the discrete-time factor model, new mathematical tools are required to deal with a continuous-time setting, where long-span asymptotics (also called increasing domain asymptotics) gives way to infill asymptotics (also called fixed domain asymptotics). For example, Fan et al. (2016) and Aït-Sahalia and Xiu (2017) extend Fan et al. (2013)'s Principal Orthogonal complement Thresholding (POET) method to high-frequency factor models. Market microstructure is an additional challenge that must be faced. The specifics of market organization and market participants' behavior induce certain short-run patterns in security prices. These patterns, such as bid-ask bounce and price-discreteness, lead to a deviation from the fundamental values (also known as efficient prices) of the securities. The security prices are thereby contaminated with market microstructure noise, which affects the estimation of parameters of interest such as volatility. Market microstructure models have been used to capture a variety of frictions inherent in the trading the efficient price process. Roll (1984) is among the first to propose a dichotomous structure in which the observed market price is the sum of the efficient price and an exogenous i.i.d. bid-ask spread. After that, Hasbrouck and Ho (1987), Choi et al. (1988) and Hasbrouck (1993) consider extended models with positive dependence in bid and ask transactions. More complicated price patterns arising from microstructure noise, such as asynchronous trading, have been investigated by researchers under the fundamental dichotomous structure.

High-dimensional models with microstructure noise have been developed more recently. Wang and Zou (2010) propose the first noise-robust estimators for the integrated volatility matrix and establish an asymptotic theory that allows both the sample size and the number of assets to approach infinity, see also Tao et al. (2011, 2013a,b), and Kim et al. (2016) for related results. However, these papers assume that the integrated volatility matrix is sparse, which often contradicts our intuition from low frequency data analysis. To solve this problem, Pelger (2019) and Dai et al. (2019) develop a continuous-time factor model with microstructure noise. Likewise, Bollerslev et al. (2019) investigate a continuous-time factor model and assume that the microstructure noise can have a factor structure itself. They use the modulated realized covariance estimator (henceforth MRC) of Christensen et al. (2010) to eliminate the effect of the microstructure noise on the estimation without explicitly estimating the factors of the microstructure noise and separating them from those of the efficient prices. They establish the consistency and bound the rate of convergence of the estimated integrated covariance matrix of the efficient price process in the large dimensional case. Related to this, Pelger (2019) classifies factors in a high-frequency factor model into jump factors and continuous factors.

We consider the dual factor model of Bollerslev et al. (2019) but we take a different approach to estimation. Our goal is to identify and separate the factors and common components from both sources: the efficient price process and the microstructure noise process. Factors for the efficient prices arise from information about future security cash flows and thereby are long-lasting, whereas factors for the microstructure noise are transient and due to the nature of trading behavior; both are of interest. We develop a methodology that is inspired by Bai and Ng (2004), who propose a test procedure called Panel Analysis of Non-stationarity in Idiosyncratic and Common Components (PANIC), which can be used to identify non-stationary factors in discrete time series. We extend the PANIC approach to our high-frequency dual factor model where the concept of co-integration cannot be used. Our methodology is in three parts. First, we estimate the common factors and loadings of both signal and noise components simultaneously from the observed return matrix. The PCA method identifies factors and idiosyncratic errors by the eigenvalues. The common factors have eigenvalues that diverge at a rate of d, where d is the number of assets, and the idiosyncratic errors have bounded eigenvalues. There is a big gap from O(d) to O(1) so that even if the efficient returns are of a smaller order, we are still able to identify their (weak) factors in the large dimensional case. The second step separates the return factors into the efficient price factors and the microstructure noise factors. This involves a second PCA step on the cumulated factors found in the first step, following the approach of Bai and Ng (2004). The final step is to cumulate the return factors to define the common factors in prices. We establish the consistency of our procedures as the number of assets increases and the number of infill observations for each asset increases. Our asymptotic framework allows for a rich diversity in the relative size of the efficient price process and the microstructure noise and in the relative size of the common component of the microstructure noise and the idiosyncratic components of the noise. This is important because a number of authors have documented that in frequently traded assets, the microstructure noise component can be quite small. Also, the Epps effect, whereby observed cross asset correlations shrink with sampling frequency, can be captured in our framework when the idiosyncratic component of the noise is larger element by element than the common component. Our model allows the so-called "weak factors", c.f., Briggs and MacCallum (2003) and Onatski (2010). We provide a full analysis of the convergence rates of all our estimators, which are affected by the magnitudes of the components of microstructure noise. To determine the number of factors we use several methods proposed in the discrete time literature and investigate their performance on simulated data. We find that the Bai and Ng (2002) information criterion performs well in terms of selecting the total number of factors and the PANIC test works well in our setting for determining the number of factors in the efficient price process. We apply our method to the intraday returns of S&P 500 Index constituents. The empirical analysis provides evidence of co-movement of the microstructure noise.

The rest of the paper is organized as follows. Section 2 specifies the model and its assumptions. Section 3 proposes the high-frequency PANIC estimation procedure and presents the asymptotic properties for the estimators. Section 3.3 provides finite-sample simulation results and Section 5 demonstrates the applicability of our proposed method through an empirical study. Section 6 concludes. The proofs of our main results are relegated to the Appendix.

Throughout the paper, we use $||\cdot||_2$ to denote the Euclidean norm of a vector. For a real symmetric matrix \mathbf{S} , we denote its k-th largest eigenvalue and trace by $\mu_k(\mathbf{S})$ and $\operatorname{tr}(\mathbf{S})$, respectively. For any $m \times n$ matrix $\mathbf{M} = (m_{ij})$, let $||\mathbf{M}||_{sp}$, $||\mathbf{M}||_1$, $||\mathbf{M}||_{\infty}$, $||\mathbf{M}||_F$ and $||\mathbf{M}||_{MAX}$ denote the spectral norm, the l_1 norm, the l_{∞} norm, the Frobenius norm, and the max norm of \mathbf{M} , respectively. Specifically, $\|\mathbf{M}\|_{sp} = \sqrt{\mu_1(\mathbf{M}^{\mathsf{T}}\mathbf{M})}$, $\||\mathbf{M}||_1 = \max_j \sum_i |m_{ij}|$, $\||\mathbf{M}||_{\infty} = \max_i \sum_j |m_{ij}|$, $\|\mathbf{M}\|_F = \sqrt{\operatorname{tr}(\mathbf{M}^{\mathsf{T}}\mathbf{M})} = \sqrt{\sum_{i,j} m_{ij}^2}$ and $\||\mathbf{M}||_{MAX} = \max_{i,j} |m_{ij}|$. Let $\mathbf{1}_n$ denote an n-dimensional vectors of 1's and \mathbf{L}_n an n-by-n lower triangular matrix where all diagonal and below-diagonal entries are 1's. Also let $a \vee b$ and $a \wedge b$ denote $\max\{a, b\}$ and $\min\{a, b\}$, and x_+ and x_- denote $\max\{0, x\}$ and $\min\{0, x\}$, respectively.

2. Model Setup and Assumptions

2.1. Dual factor structure

Let X_{it} denote the observed log transaction price of stock *i* at time *t*, for i = 1, ..., d. We allow *d* to diverge with *n*, although we have suppressed the subscript *n* for *d*. For the sake of simplicity, we assume that price observations of all stocks are synchronously collected, and that price observations for each stock are equidistantly collected in the fixed time interval [0, T]. Thus we do not consider non-synchronous trading explicitly. Without loss of generality, we let T = 1. Let *n* be the number of observations and $\Delta = 1/n$. Then, the prices are observed at the time points $t = 0, \Delta, 2\Delta, \ldots, n\Delta$. Although we assume the length of the time interval between any two successive transactions to be equal, our results still hold under the more general assumption that the length of the time interval is uniformly of order 1/n.

We assume that the observed log transaction price, X_{it} , can be decomposed into the unobserved efficient log-price X_{it}^* plus a noise component Z_{it} , i.e.,

$$X_{it} = X_{it}^* + Z_{it} \quad \text{or} \quad \boldsymbol{X}_t = \boldsymbol{X}_t^* + \boldsymbol{Z}_t, \tag{2.1}$$

where $\mathbf{X}_{t}^{*} = (X_{1t}^{*}, \ldots, X_{dt}^{*})^{\mathsf{T}}$ and $\mathbf{Z}_{t} = (Z_{1t}, \ldots, Z_{dt})^{\mathsf{T}}$. For each component of X_{it} , we introduce a factor structure (see Assumptions 1 and 2 below) and therefore, name the model as a *dual factor* model.

Assumption 1. (Factor Structure for Efficient Log-price) (i) The efficient log-price X_t^* follows a factor model of the form,

$$\left\{ egin{array}{rcl} {\mathsf d} {oldsymbol X}^*_t &=& {oldsymbol \Lambda}_F {\mathsf d} {oldsymbol F}_t + {\mathsf d} {oldsymbol U}_t, \ {\mathsf d} {oldsymbol F}_t &=& {oldsymbol \sigma}_{Ft} {\mathsf d} {oldsymbol B}^U_t, \ {\mathsf d} {oldsymbol U}_t &=& {oldsymbol \sigma}_{Ut} {\mathsf d} {oldsymbol B}^U_t, \end{array}
ight.$$

where $\mathbf{\Lambda}_F = (\lambda_{F,ik})_{1 \leq i \leq d, 1 \leq k \leq K_F}$ denotes the $d \times K_F$ matrix of factor loadings, K_F is the number of factors, $\mathbf{F}_t = (F_{1t}, ..., F_{K_Ft})^{\mathsf{T}}$ denotes latent factors, $\mathbf{U}_t = (U_{1t}, ..., U_{dt})^{\mathsf{T}}$ is the idiosyncratic component, $\boldsymbol{\sigma}_{Ft}$ is a $K_F \times K_F$ càd-làg spot volatility matrix for factors, $\boldsymbol{\sigma}_{Ut}$ is a $d \times d$ càd-làg spot volatility matrix for idiosyncratic errors, and $\mathbf{B}_t^F = (B_{1t}^F, ..., B_{K_Ft}^F)^{\mathsf{T}}$ and $\mathbf{B}_t^U = (B_{1t}^U, ..., B_{dt}^U)^{\mathsf{T}}$ are independent Brownian motions.

(ii) The initial states of \mathbf{F}_t and \mathbf{U}_t satisfy $\|\mathbf{F}_0\|_{MAX} = O_P(1)$ and $\|\mathbf{U}_0\|_{MAX} = O_P(1)$.

(iii) Denote the spot covariance matrices of \mathbf{F}_t and \mathbf{U}_t by $\Sigma_{Ft} = \boldsymbol{\sigma}_{Ft}\boldsymbol{\sigma}_{Ft}^{\mathsf{T}}$ and $\Sigma_{Ut} = \boldsymbol{\sigma}_{Ut}\boldsymbol{\sigma}_{Ut}^{\mathsf{T}}$, and define $\Sigma_{Ft-} = \lim_{\Delta \to 0} \Sigma_{F,t-\Delta}$ and $\Sigma_{Ut-} = \lim_{\Delta \to 0} \Sigma_{U,t-\Delta}$. Then, Σ_{Ft} , Σ_{Ft-} , Σ_{Ut} and Σ_{Ut-} are all positive-definite and satisfy

 $\max\{\|\boldsymbol{\Sigma}_{Ft}\|_{\mathsf{MAX}}, \|\boldsymbol{\Sigma}_{Ft-}\|_{\mathsf{MAX}}, \|\boldsymbol{\Sigma}_{Ut}\|_{\mathsf{MAX}}, \|\boldsymbol{\Sigma}_{Ut-}\|_{\mathsf{MAX}}\} \le C_{\sigma}$

for all $t \in [0,1]$, where C_{σ} is a positive constant independent of n and d.

Remark 2.1. Our dual factor model follows the setting of Bollerslev et al. (2019), and inherits several limitations. Firstly, we do not allow a drift term in the diffusion model. More general settings can be seen in Dai et al. (2019) and Barigozzi et al. (2020b) for instance. Secondly, we assume that the factor loadings are constant and thus, neither random nor time-varying loadings are allowed. We refer to Aït-Sahalia et al. (2020) for high-frequency factor models with time-varying betas. Thirdly, we assume the number of factors is fixed. We refer to Fan et al. (2011, 2016) for factor models with the number of factors increasing with d. Another limitation is that jumps are excluded in our model. Dealing with jumps would require more complicated procedures and hence, we leave it for future work.

Assumption 2. (Factor Structure for Market Microstructure Noise)

The microstructure noise Z_t follows a factor model whose magnitude may depend on the sampling frequency, that is

$$\boldsymbol{Z}_t = \boldsymbol{\Lambda}_G \boldsymbol{D}_G \boldsymbol{G}_t + \boldsymbol{D}_V \boldsymbol{V}_t, \qquad (2.2)$$

where $\mathbf{\Lambda}_G = (\lambda_{G,ik})_{1 \leq i \leq d,1 \leq k \leq K_G}$ denotes the $d \times K_G$ matrix of factor loadings with K_G being the number of factors for microstructure noise, $\mathbf{G}_t = (G_{1t}, ..., G_{K_Gt})^{\mathsf{T}}$ denotes the latent factors, $\mathbf{V}_t = (V_{1t}, ..., V_{dt})^{\mathsf{T}}$ is the vector of idiosyncratic components, and \mathbf{D}_G and \mathbf{D}_V are two diagonal matrices satisfying $\mu_1(\mathbf{D}_G) = O(n^{\bar{\tau}_G})$, $\mu_1(\mathbf{D}_G^{-1}) = O(n^{-\underline{\tau}_G})$ and $\mu_1(\mathbf{D}_V) = O(n^{\bar{\tau}_V})$, where $\bar{\tau}_G$, $\underline{\tau}_G$, and $\bar{\tau}_V$ are some constants.

The introduction of \mathbf{D}_G and \mathbf{D}_V in (2.2) allows the microstructure noise to be larger or smaller in magnitude than the efficient log-prices, depending on the values of $\bar{\tau}_G$ and $\bar{\tau}_V$ (see also Kalnina and Linton (2008) for a model where the microstructure noise can be large or small, depending on the value of a magnitude parameter). With such a setup, the magnitude of $\mathbf{\Lambda}_G$, \mathbf{G}_t and \mathbf{V}_t is independent of n. A similar treatment can be found in Kim et al. (2016). Our model is an extension of the model of Bollerslev et al. (2019), who only consider $\mathbf{D}_G = \mathbf{I}_{K_G}$ and $\mathbf{D}_V = \mathbf{I}_d$. However, they use a similar setting when generating simulation data (also see Section 3.3) without discussing the asymptotic impacts of \mathbf{D}_G and \mathbf{D}_V .

To introduce the first-differenced form of the dual factor model, we use little letters to denote the first-order differences of random variables. Specifically, define the return as $x_t = X_t - X_{t-\Delta}$, and the

efficient return (or frictionless return) as $\boldsymbol{x}_{t}^{*} = \int_{t-\Delta}^{t} d\boldsymbol{X}_{s}^{*} = \boldsymbol{X}_{t}^{*} - \boldsymbol{X}_{t-\Delta}^{*}$. Denote $\boldsymbol{f}_{t} = \int_{t-\Delta}^{t} \boldsymbol{\sigma}_{fs} d\boldsymbol{B}_{s}^{F} = \boldsymbol{F}_{t} - \boldsymbol{F}_{t-\Delta}, \, \boldsymbol{g}_{t} = \boldsymbol{G}_{t} - \boldsymbol{G}_{t-\Delta}, \, \boldsymbol{z}_{t} = \boldsymbol{Z}_{t} - \boldsymbol{Z}_{t-\Delta}, \, \boldsymbol{u}_{t} = \int_{t-\Delta}^{t} \boldsymbol{\sigma}_{Us} d\boldsymbol{B}_{s}^{U} = \boldsymbol{U}_{t} - \boldsymbol{U}_{t-\Delta}, \, \text{and} \, \boldsymbol{v}_{t} = \boldsymbol{V}_{t} - \boldsymbol{V}_{t-\Delta}.$ Then by (2.1)–(2.2), we can write the dual factor model as

$$\begin{cases} \boldsymbol{x}_t = \boldsymbol{x}_t^* + \boldsymbol{z}_t \\ \boldsymbol{x}_t^* = \boldsymbol{\Lambda}_F \boldsymbol{f}_t + \boldsymbol{u}_t \\ \boldsymbol{z}_t = \boldsymbol{\Lambda}_G \boldsymbol{D}_G \boldsymbol{g}_t + \boldsymbol{D}_V \boldsymbol{v}_t \end{cases}$$
(2.3)

Combining the factor structures for \boldsymbol{x}_t^* and \boldsymbol{z}_t , we have

$$\boldsymbol{x}_{t} = \boldsymbol{\Lambda}_{F} \boldsymbol{f}_{t} + \boldsymbol{u}_{t} + \boldsymbol{\Lambda}_{G} \boldsymbol{D}_{G} \boldsymbol{g}_{t} + \boldsymbol{D}_{V} \boldsymbol{v}_{t},$$

= $\boldsymbol{\Lambda}_{H} \boldsymbol{D}_{H} \boldsymbol{h}_{t} + \boldsymbol{w}_{t},$ (2.4)

where $\mathbf{h}_t = (\mathbf{f}_t^{\mathsf{T}}, n^{-1/2} \mathbf{g}_t^{\mathsf{T}})^{\mathsf{T}}$, $\mathbf{\Lambda}_H = (\mathbf{\Lambda}_F, \mathbf{\Lambda}_G)$, $\mathbf{w}_t = \mathbf{u}_t + \mathbf{D}_V \mathbf{v}_t$ and $\mathbf{D}_H = \mathsf{diag}(\mathbf{I}_{K_F}, n^{1/2} \mathbf{D}_G)$. This can be seen as a factor structure for \mathbf{x}_t with $\mathbf{h}_t = (\mathbf{f}_t^{\mathsf{T}}, n^{-1/2} \mathbf{g}_t^{\mathsf{T}})^{\mathsf{T}}$ being the factors and $\mathbf{\Lambda}_H = (\mathbf{\Lambda}_F, \mathbf{\Lambda}_G)$ being the factor loadings. Note that \mathbf{g}_t is divided by $n^{1/2}$ in \mathbf{h}_t so that both components of \mathbf{h}_t are of the same magnitude. Consequently, the magnitude matrix \mathbf{D}_G is multiplied by $n^{1/2}$ in \mathbf{D}_H . This also gives more insight into the role \mathbf{D}_G plays. If $\bar{\tau}_G > -1/2$, the factors of the microstructure noise dominate those of the efficient returns and vice versa. Similarly, since \mathbf{u}_t and $n^{-1/2}\mathbf{v}_t$ are of the same magnitude, if $\bar{\tau}_V > -1/2$, the idiosyncratic components of the microstructure noise dominates those of the efficient returns and vice versa.

Denote the number of independent factors in factor model (2.4) as K_H . If f_t and g_t are collinear, K_H will be less than $K_F + K_G$. In such a case, the efficient prices and microstructure noise are not separable and the dichotomous structure fails. Thus for identification purposes, we exclude this situation and make the following assumption.

Assumption 3. (Independence between efficient prices and microstructure noise) The discrete time series G_t and V_t are independent of the continuous-time processes F_t and U_t .

Assumption 4. (Factor loadings) The factor loadings matrix Λ_H satisfies

 $\|\mathbf{\Lambda}_H\|_{\mathsf{MAX}} < C_{\Lambda}, \qquad \|\mathbf{\Lambda}_H^{^{\mathsf{T}}}\mathbf{\Lambda}_H/d\|_{\mathsf{sp}} < C_{\Lambda}, \qquad \|(\mathbf{\Lambda}_H^{^{\mathsf{T}}}\mathbf{\Lambda}_H/d)^{-1}\|_{\mathsf{sp}} < C_{\Lambda},$

where C_{Λ} is a positive constant independent of n and d.

Remark 2.2. If σ_{Ft} is a constant spot volatility matrix that does not vary with t, then we can use the notion of cointegration as in Bai and Ng (2004) when $T \to \infty$. In this case, the cointegration rank of $H_t := (\mathbf{F}_t^{\mathsf{T}}, n^{-1/2} \mathbf{G}_t^{\mathsf{T}})^{\mathsf{T}}$ is $K_H - K_F = K_G$, which is the number of factors for microstructure noise.

Remark 2.3. It is prevalent to assume independence between price components due to fundamental security value and noise attributable to market rules and trading mechanisms in a univariate microstructure model. One of the reasons is modelling simplicity, so that the two components of the model can be identified and the model can have sensible economic and statistical representation and interpretation. However, the independence assumption may sometimes be violated. For example,

it is shown in Delattre and Jacod (1997) and Li and Mykland (2007) that this assumption can be substantially weakened, especially for processes contaminated with rounding errors. Glosten (1987) and Glosten and Harris (1988) also show that the microstructure noise may no longer be uncorrelated with the efficient price when asymmetric information is involved. Recently, there has been some research that allows correlation between them, such as Kunitomo and Kurisu (2019), who modify the Separating Information Maximum Likelihood due to Kunitomo and Sato (2013) to detect hidden factors of quadratic variation in the presence of correlated market microstructure noise.

For easy reference, we summarise the notation used for variables and factors in Table 1. We use different fonts to distinguish between matrices, vectors and scalars. For example, $\mathscr{X} = (X_{\Delta}, \ldots, X_{n\Delta})^{\mathsf{T}}$ is an $n \times d$ matrix of observed prices, $X_{s\Delta}^{\mathsf{T}}$ is its s-th row, and $X_{i,s\Delta}$ is the (s, i)-entry of the matrix \mathscr{X} . Following the same rule, other variables are defined analogously.

Variables	Aggr	regation	Form	First-o	differen	ce Form
	Matrix	Row	Element	Matrix	Row	Element
	-wise	-wise	-wise	-wise	-wise	-wise
Observed price	\mathscr{X}	$oldsymbol{X}_t$	X_{it}	x	$oldsymbol{x}_t$	x_{it}
Efficient price (EP)	\mathscr{X}^*	$oldsymbol{X}_t^*$	X_{it}^*	x^*	$oldsymbol{x}_t^*$	x_{it}^*
Microstructure noise (MN)	\mathscr{Z}	$oldsymbol{Z}_t$	Z_{it}	z	$oldsymbol{z}_t$	z_{it}
Factors for EP	${\cal F}$	$oldsymbol{F}_t$	F_{jt}	f	$oldsymbol{f}_t$	f_{jt}
Factors for MN	${\mathcal G}$	$oldsymbol{G}_t$	G_{jt}	q	$oldsymbol{g}_t$	g_{jt}
Total factors	${\mathcal H}$	$oldsymbol{H}_t$	$\dot{H_{it}}$	ħ	$oldsymbol{h}_t$	$\dot{h_{it}}$
Idiosyncratic errors for EP	${\mathscr U}$	$oldsymbol{U}_t$	U_{it}	U	$oldsymbol{u}_t$	$\dot{u_{it}}$
Idiosyncratic errors for MN	V	$oldsymbol{V}_t$	V_{it}	v	$oldsymbol{v}_t$	v_{it}
Total idiosyncratic errors	W	$oldsymbol{W}_t$	W_{it}	W	$oldsymbol{w}_t$	w_{it}

Table 1: Notation for variables/factors in the dual factor model

¹ The first dimension of the matrices in the table are set as time, while the second dimension are set as a stock or a factor. We have $t = \Delta, \dots, n\Delta, 1 \le i \le d$ and $1 \le j \le K$, where $K = K_F$, K_G , or K_H , depending on the circumstance. Note that in the subscript of the element-wise notation, we write the column index first.

² Although when t = 0, X_t is observable, t starts from Δ in the matrix-wise notation for both aggregation form and first-difference form, for the sake of consistency.

2.2. Covariance structure

Define the integrated covolatility matrix of f_t and u_t as Σ_F and Σ_U , respectively. That is $\Sigma_F = \int_0^1 \Sigma_{Ft} dt$ and $\Sigma_U = \int_0^1 \Sigma_{Ut} dt$. Then the factor structure for efficient prices leads to the following identity,

$$\Sigma_{x^*} = \Lambda_F \Sigma_F \Lambda_F^{\mathsf{T}} + \Sigma_U, \qquad (2.5)$$

where Σ_{x^*} is the integrated covolatility matrix of x_t^* . Similarly, the factor structure for microstructure noise leads to the following identity,

$$\Sigma_z = \Lambda_G \mathbf{D}_G \Sigma_g \mathbf{D}_G \Lambda_G^{\mathsf{T}} + \mathbf{D}_V \Sigma_v \mathbf{D}_V, \qquad (2.6)$$

where $\Sigma_z = \text{Var}(\boldsymbol{z}_t)$, $\Sigma_g = \text{Var}(\boldsymbol{g}_t)$, and $\Sigma_v = \text{Var}(\boldsymbol{v}_t)$.¹ For the integrated covolatility matrix of observed prices, we combine (2.5) and (2.6) to obtain

$$\boldsymbol{\Sigma}_{x} = \boldsymbol{\Lambda}_{H} \mathbf{D}_{H} \boldsymbol{\Sigma}_{h} \mathbf{D}_{H} \boldsymbol{\Lambda}_{H}^{\mathsf{T}} + \boldsymbol{\Sigma}_{w}, \qquad (2.7)$$

where $\Sigma_x = \Sigma_{x^*} + n\Sigma_z$, $\Sigma_h = \text{diag}(\Sigma_F, \Sigma_g)$, and $\Sigma_w = \Sigma_U + n\mathbf{D}_V \Sigma_v \mathbf{D}_V$. The reason we multiply Σ_z by *n* is to be consistent with its sample estimates. Note that $\widehat{\Sigma}_z = n^{-1} \sum_{i=1}^n \mathbf{z}_t \mathbf{z}_t^{\mathsf{T}}$ is an approximation of Σ_z while $\widehat{\Sigma}_{x^*} = \sum_{i=1}^n (\mathbf{x}_t^*) (\mathbf{x}_t^*)^{\mathsf{T}}$ is an approximation of Σ_{x^*} .

In order to identify the factor structure, we make the following sparsity assumption, as in Fan et al. (2013) and Aït-Sahalia and Xiu (2017). For ease of composition, we define

$$\bar{\tau}_G^{*+} = (1/2 + \bar{\tau}_G)_+, \quad \underline{\tau}_G^{*-} = (1/2 + \underline{\tau}_G)_-, \quad \bar{\tau}_V^* = (1/2 + \bar{\tau}_V), \quad \text{and } \bar{\tau}_V^{*+} = (1/2 + \bar{\tau}_V)_+.$$

Assumption 5. (Sparsity of idiosyncratic integrated covariance matrices) The integrated volatility matrices of the idiosyncratic components satisfy

$$\|\boldsymbol{\Sigma}_{U}\|_{1} = O(m_{U,d}), \quad \|\boldsymbol{\Sigma}_{v}\|_{1} = O(m_{v,d}), \quad m_{w,nd}/(d^{1/2}n^{2\underline{\tau}_{G}^{*-}}) \to 0 \text{ with } m_{w,nd} = m_{U,d} + n^{2\bar{\tau}_{V}^{*}}m_{v,d}.$$

Under Assumption 5, we have $\|\Sigma_w\|_1 = O(m_{w,nd})$. It is worth pointing out that the condition $m_{w,nd}/(dn^{2\underline{\tau}_G^{*-}}) \to 0$ is sufficient for the identification of the factors, see Lemma A.1. When only the weaker condition is required, we refer to the assumption as Assumption 5^{*}. In order to obtain consistent estimates, we further require $m_{w,nd}/(d^{1/2}n^{2\underline{\tau}_G^{*-}}) \to 0$ as in Assumption 5. That estimation of an approximate factor model requires more strict sparsity conditions than identification is also observed in Aït-Sahalia and Xiu (2017).

We also make the following assumption on the stationarity of the microstructure noise components.

Assumption 6. (Stationary and sub-Gaussian microstructure noise) (i) The series $\{\mathbf{G}_t, \mathbf{V}_t\}$ is strictly stationary. In addition, $E[G_{jt}] = E[V_{it}] = E[G_{jt}V_{it}] = 0$ for all $1 \le i \le d, 1 \le j \le K_G$ and $t = 0, \Delta, \dots, n\Delta$.

(ii) There exist positive constants $C_{\alpha} > 0$ and $\gamma_1 > 0$ such that the strong mixing sequence $\alpha(.)$ of the series $\{\mathbf{G}_t, \mathbf{V}_t\}$ satisfies $\alpha(s\Delta) \leq C_{\alpha} \exp(-s^{\gamma_1})$.

(iii) There exist $b_1 > 0$, $\gamma_2 > 0$, with $\gamma_1^{-1} + 3\gamma_2^{-1} > 1$, such that for all c > 0, we have

$$\max_{1 \le j \le K_G} P\left(|G_{jt}| > c\right) \le \exp(1 - (c/b_1)^{\gamma_2})$$
(2.8)

and

$$\max_{1 \le i \le d} P(|V_{it}| > c) \le \exp(1 - (c/b_1)^{\gamma_2}),$$
(2.9)

¹The rules of using capital letters in the subscript are that (i) For factor loading matrices, magnitude matrices and number of factors, we use capital letters; (ii) For the integral volatility matrices of a continuouse processes, we use capital letters; (iii) For other cases such as the covariance matrices of discrete time series and the contaminated integral volatility matrices, e.g., Σ_{x^*} and Σ_z , we choose according to the notation of the series.

for $t = 0, \Delta, \ldots, n\Delta$. (iv) There exist $b_2 > 0$, $\gamma_3 > 0$, with $\gamma_1^{-1} + 3\gamma_3^{-1} > 1$, such that for all c > 0, we have

$$\max_{1 \le j \le K_H} P\left(d^{-1/2} |\boldsymbol{\lambda}_{H,\cdot,j}^{\mathsf{T}} \boldsymbol{V}_t| > m_{v,d}^{1/2} \cdot c\right) \le \exp(1 - (c/b_2)^{\gamma_3}),\tag{2.10}$$

for $t = 0, \Delta, ..., n\Delta$, where $\lambda_{H,.,j}$ is the *j*-th column of Λ_H ; (v) Let $1/\gamma = 1/\gamma_1 + 3/(\gamma_2 \vee \gamma_3)$. Then, $(\log d)^{2/\gamma - 1} = o(n)$.

Assumptions 6(i) and (ii) are similar to Assumption 3.2(i) and Assumption 3.3 in Fan et al. (2013). The strong mixing condition is more general than Assumptions 1 and 3 in Barigozzi et al. (2020b) in which the common components and idiosyncratic components of the microstructure noise are both assumed to be linear processes. The series G_t and V_t can be serially correlated, which is more general than the assumption in Kim et al. (2016). Moreover, V_t can be cross-sectionally dependent, in which case, (2.2) gives an approximate factor model for the microstructure noise. More general assumptions can be found in Bai and Ng (2002) which permit weakly correlated idiosyncratic errors. Assumption 6(iii) requires exponential-type tails for the distributions of G_t and V_t , which is similar to Assumption 3.2(iii) of Fan et al. (2013) and Condition A1 of Tao et al. (2013b). If the microstructure noise has a heavier tail and the sub-Gaussian assumption is violated, robust estimation can be used to mitigate the influence of heavy tails (see, for example Fan and Kim (2018)). But to focus on the key objective of this paper (i.e., to identify and estimate factors for both efficient price and microstructure noise), we do not consider it in this paper.

Assumption 6(iv) is an additional exponential tail condition. Intuitively, let us consider two extreme cases. If v_t is cross-sectionally jointly independent, (2.10) holds true with $m_{v,d} = 1$. If v_t is cross-sectionally comonotonic, (2.10) holds true with $m_{v,d} = d$. The sparsity condition in Assumption 5 rules out the second case. But v_t can be cross-sectionally weakly dependent, with $m_{v,d}$ between 1 and $d^{1/2}n^{-1-2\tau_V+(1+2\tau_G)-}$. Assumption 6(iv) guarantees that

$$\left\| (nd)^{-1} \sum_{s=1}^{n} \mathbf{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{v}_{s\Delta} \boldsymbol{v}_{s\Delta}^{\mathsf{T}} \mathbf{\Lambda}_{H} - d^{-1} \mathbf{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{\Sigma}_{v} \mathbf{\Lambda}_{H} \right\|_{\mathsf{MAX}} = O_{P}(m_{v,d} (\log d/n)^{1/2}),$$

see Lemma B.2(iv). Assumption 6(iv) can also be seen as an extension of the condition (see Assumption F.3 of Bai (2003) and Assumption 3.4(iii) of Fan et al. (2013))

$$\max_{1 \le j \le K_H} \mathsf{E}[(d^{-1/2}|\boldsymbol{\lambda}_{H,\cdot,j}^{\mathsf{T}}\boldsymbol{v}_t|)^4] < C,$$

for some positive constant C. The latter condition is no longer true when $d^{-1} \| \mathbf{\Lambda}_H^{\mathsf{T}} \mathbf{\Sigma}_v \mathbf{\Lambda}_H \|_{\mathsf{MAX}} \to \infty$. Assumption 6(v) presents a trade-off between the mixing and tail conditions and the dimension d.

3. Estimation Procedure and Asymptotic Results

In this section, we develop a three-step estimation procedure for the common factors, \mathscr{F} and \mathscr{G} , of the dual factor model. The estimation procedure is based on two PCA procedures and so we call it Double PCA or DPCA. The asymptotic results are presented step by step, so as to provide

a better understanding on what the intermediate estimators estimate. In the first step, the factors and factor loadings for the combined factor model (2.4) in the first-difference form are estimated. In the second step, we cumulate the factors and separate the factors of the efficient prices from those of the microstructure noise in the first-difference form. In the last step, we cumulate \hat{f}_t and \hat{g}_t to get \hat{F}_t and \hat{G}_t , respectively. For the time being, we assume that the number of factors K_F and K_G are known, and $K_H = K_F + K_G$. We will discuss how to determine them in Section 4.1.

An additional technical condition is required, which is similar to Assumption G in Bai (2003) and Assumption 4 in Aït-Sahalia and Xiu (2017) in the case of $\mathbf{D}_H = \mathbf{I}_{K_H}$. Note that we have assumed the boundedness of the largest eigenvalue of the matrix $d^{-1} \mathbf{\Lambda}_H^{\mathsf{T}} \mathbf{\Lambda}_H$ in Assumption 4, instead of the convergence of the matrix. Hence, we do not introduce a limiting matrix for it.

Assumption 7. (Distinct eigenvalues)

The $d \times d$ matrix $\Lambda_H \mathbf{D}_H \boldsymbol{\Sigma}_h \mathbf{D}_H \Lambda_H^{\mathsf{T}}$ has asymptotically distinct eigenvalues in the sense that for $j = 2, \ldots d$

$$\frac{\mu_{j-1}(\boldsymbol{\Lambda}_{H}\boldsymbol{\mathrm{D}}_{H}\boldsymbol{\Sigma}_{h}\boldsymbol{\mathrm{D}}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}) - \mu_{j}(\boldsymbol{\Lambda}_{H}\boldsymbol{\mathrm{D}}_{H}\boldsymbol{\Sigma}_{h}\boldsymbol{\mathrm{D}}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}})}{\mu_{K_{H}}(\boldsymbol{\Lambda}_{H}\boldsymbol{\mathrm{D}}_{H}\boldsymbol{\Sigma}_{h}\boldsymbol{\mathrm{D}}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}})} > C_{2}$$

when $(n, d) \to \infty$, for $j = 2, ..., K_H$, where C is a positive constant independent of n and d.

We also require the following assumption, which restricts the relation between n and d.

Assumption 8. (Relation between n and d) As $n \to \infty$, (i)

$$n^{1+4\underline{\tau}_G^{*-}-4(\bar{\tau}_G^{*+}\vee\bar{\tau}_V^{*+})}/(\log d)\to\infty,$$

and

(ii)

$$n^{1+4\underline{\tau}_{G}^{*-}-4(\bar{\tau}_{G}^{*+}\vee\bar{\tau}_{V}^{*+})}m_{w,nd}^{2}/(d^{2}\log d) \to 0.$$

When $\bar{\tau}_G = \underline{\tau}_G = \bar{\tau}_V = -1/2$ and $m_{w,nd} = O(1)$, Assumption 8 degenerates to $n/(\log d) \to \infty$ and $n/(d^2 \log d) \to 0$, which are similar to the corresponding condition in Theorem 3.1 of Fan et al. (2013) and assumption A1 of Tao et al. (2013b).

3.1. First step: PCA estimation in first-difference form

The first step in our estimation procedure is to apply PCA to the first-difference form of the dual factor model to extract estimates of all the factors and factor loadings (for both efficient prices and microstructure noise).

Recall that $x = \hbar \Lambda_H^{\mathsf{T}} + w$. Different identification conditions for the factor structure can be used, according to whether we would like to normalise the factors or the factor loadings. Accordingly, there are two ways to implement a PCA to estimate the factors and factor loadings. If we normalise the factor loadings, the PCA estimator solves the following optimisation problem,

$$\begin{cases} \min_{\boldsymbol{\Lambda}_{H,\mathcal{R}}} & \|\boldsymbol{x} - \boldsymbol{\mathcal{R}} \boldsymbol{\Lambda}_{H}^{\mathsf{T}} \|_{F}, \\ \text{s.t.} & \boldsymbol{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{\Lambda}_{H} / d = \mathbf{I}_{K_{H}}. \end{cases}$$
(3.1)

Computationally, we conduct an eigen-decomposition of the $d \times d$ matrix $x^{\mathsf{T}}x$, and obtain $\widehat{\Lambda}_H = (\widehat{\lambda}_{H1}, \ldots, \widehat{\lambda}_{Hd})^{\mathsf{T}}$, which is the $d \times K_H$ matrix consisting of the K_H eigenvectors (multiplied by \sqrt{d}) corresponding to the K_H largest eigenvalues of $x^{\mathsf{T}}x$. The common factors can be estimated as

$$\widehat{\hbar} = x\widehat{\Lambda}_H/d = \left(\widehat{h}_{\Delta}, \dots, \widehat{h}_{n\Delta}\right)^{\mathsf{T}}.$$
(3.2)

Alternatively, if we normalise the factors, the PCA estimator solves the following optimisation problem,

$$\left\{ egin{array}{cc} \min & \|x- \hbar {f \Lambda}_H^{^{\intercal}}\|_F, \ {f s.t.} & \hbar^{^{\intercal}} \hbar = {f I}_{K_H}. \end{array}
ight.$$

Computationally, we conduct PCA^{*} (we add an asterisk to distinguish it from the PCA above that is based on the normalisation of the factor loadings) on the $n \times n$ matrix xx^{T} , and get $\hat{k}^* = (\hat{h}^*_{\Delta}, \ldots, \hat{h}^*_{n\Delta})^{\mathsf{T}}$ denoting the $n \times K_H$ matrix consisting of the K_H eigenvectors corresponding to the K_H largest eigenvalues of xx^{T} . The common factors can be estimated as

$$\widehat{\Lambda}_{H}^{*} = x^{\mathsf{T}} \widehat{\mathscr{R}}^{*} = \left(\widehat{\lambda}_{H1}^{*}, \dots, \widehat{\lambda}_{Hd}^{*}\right)^{\mathsf{T}}.$$
(3.3)

It is easy to show that both \hat{k} and \hat{k}^* are eigenvectors of xx^{T} , and thereby, they span the same space. In this sense, the two ways of PCA are equivalent. Bai and Li (2012) point out that the analysis of one PCA representation will carry over to the other by switching the roles of n and d and the role of the factor loadings and factors, and thus it is sufficient to carefully examine the asymptotic properties of one representation. All high-frequency factor models, as we know, use (3.1) to implement PCA. The main reason may be that the spot covariance matrices is time-varying and thereby the $n \times n$ matrix xx^{T} is conceptually more difficult to analyse.

However, in our three-step estimation procedure, the final estimators based on the above two PCAs will not be equivalent. Since the first-step factor estimators will be fed into the second step in cumulative form for another PCA, whether the factors are normalised in the first step will affect the final results. In Section 3.3, we will compare the small sample performance of estimators based on the two different PCA representations in the first step. We will see that PCA* always outperforms PCA. Hence, we will use PCA* for our first step and derive the asymptotic theory based on it.

Since the asymptotic theory of $\hat{\mathscr{R}}$ is easier to establish than that of $\hat{\mathscr{R}}^*$, as discussed above, we first prove the consistency of $\hat{\mathscr{R}}$ (see Lemma A.2 in Appendix) and then use the following relation

$$\widehat{\hbar}^* = \widehat{\hbar} (\widehat{\hbar}^{\mathsf{T}} \widehat{\hbar})^{-1/2} \quad \text{and} \quad \widehat{\Lambda}^*_H = \widehat{\Lambda}_H (\widehat{\hbar}^{\mathsf{T}} \widehat{\hbar})^{1/2}, \tag{3.4}$$

to prove the consistency of $\hat{\mathscr{R}}^*$. When *d* is much larger than *n*, it is computationally more convenient to conduct PCA on the $n \times n$ matrix $\mathscr{x} \mathscr{x}^{\mathsf{T}}$, and vice versa. Then the relations in (3.4) can be used to get the desired form of estimates.

The following theorem shows the uniform rate of convergence for $\widehat{\Lambda}_{H}^{*}$ and $\widehat{\mathscr{K}}^{*}$ of the dual factor model. For ease of composition, we denote

$$a_{nd} = (\log d)^{1/2} \frac{n^{\bar{\tau}_V^{*+} + \bar{\tau}_G^{*+} \vee \bar{\tau}_V^{*+}}}{n^{1/2}} + \frac{m_{w,nd}}{d^{1/2}}.$$
(3.5)

Theorem 3.1. Suppose that Assumptions 1-8 are satisfied. We have (i)

$$\left\|\widehat{\mathscr{R}}^{*^{\mathsf{T}}} - (\mathbf{R}^{*})^{-1}(\mathscr{R}^{\mathsf{T}}\mathscr{R})^{-1/2}\mathscr{R}^{\mathsf{T}}\right\|_{\mathsf{MAX}} \leq \left\|\widehat{\mathscr{R}}^{*^{\mathsf{T}}} - (\mathbf{R}^{*})^{-1}(\mathscr{R}^{\mathsf{T}}\mathscr{R})^{-1/2}\mathscr{R}^{\mathsf{T}}\right\|_{\mathsf{sp}} = O_P\left(n^{-2\underline{\tau}_{G}^{*^{-}}} \cdot a_{nd}\right), \quad (3.6)$$

where the rotation matrix \mathbf{R}^* is defined by

$$\mathbf{R}^* = d^{1/2} (\boldsymbol{\hbar}^{\mathsf{T}} \boldsymbol{\hbar})^{-1/2} \boldsymbol{\Lambda}_H^{\mathsf{T}} \widehat{\boldsymbol{\Lambda}}_H \widehat{\mathbf{D}}_{x, K_H}^{-1/2}$$

in which $\widehat{\mathbf{D}}_{x,K_H} = d\widehat{\lambda}^{\mathsf{T}} \widehat{\lambda} = \widehat{\Lambda}_H^{\mathsf{T}} \widehat{\Lambda}_H^*$ is a $K_H \times K_H$ diagonal matrix with the diagonal elements being the first K_H largest eigenvalues of $x^{\mathsf{T}} x$ arranged in a descending order; (ii)

$$\left\|\widehat{\mathbf{\Lambda}}_{H}^{*}-\mathbf{\Lambda}_{H}\mathbf{D}_{H}(\boldsymbol{\mathscr{R}}^{\mathsf{T}}\boldsymbol{\mathscr{R}})^{1/2}\mathbf{R}^{*}\right\|_{\mathsf{MAX}}=O_{P}\left(n^{-2\underline{\tau}_{G}^{*-}}\cdot a_{nd}\right);$$
(3.7)

(iii) $\widehat{\lambda}^* \widehat{\Lambda}_H^{*^{\mathsf{T}}} = \widehat{\lambda} \widehat{\Lambda}_H^{\mathsf{T}}$ and

$$\left\|\widehat{\boldsymbol{\lambda}}^{*}\widehat{\boldsymbol{\Lambda}}_{H}^{*^{\mathsf{T}}} - \boldsymbol{\boldsymbol{\hbar}}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\right\|_{\mathsf{MAX}} = O_{P}\left(\boldsymbol{n}^{\overline{\tau}_{G}^{*^{+}} - 2\underline{\tau}_{G}^{*^{-}}} \cdot \boldsymbol{a}_{nd}\right);$$
(3.8)

(iv) \mathbf{R}^* is an asymptotically orthogonal matrix, that is

$$\|\mathbf{R}^{*^{\mathsf{T}}}\mathbf{R}^{*}-\mathbf{I}_{K_{H}}\|_{\mathsf{sp}}=O_{P}(n^{-2\underline{\tau}_{G}^{*-}}\cdot a_{nd}).$$

Note that $n^{-2\underline{\tau}_{G}^{*-}}a_{nd} = o(1)$ under Assumption 8. The theorem shows the uniform rate of convergence for $\hat{\mathbb{A}}^{*}$ of the dual factor model, which is similar to Theorem 5 in Aït-Sahalia and Xiu (2017) for high-frequency factor model without the magnitude matrices. The estimator $\hat{\mathbb{A}}^{*}$ converges to normalised factors $(\mathbb{A}^{\dagger}\mathbb{A})^{-1/2}\mathbb{A}^{\dagger}$ up to an asymptotically orthogonal matrix. The result is different from the uniform results for low-frequency factor models, e.g., Proposition 2 in Bai (2003) and Theorem 3.3 in Fan et al. (2013). In a low-frequency factor model, the uniform convergence of the estimator for $n^{1/2}(\mathbb{A}^{\dagger}\mathbb{A})^{-1/2}\mathbb{A}^{\dagger}$ can be derived, while in a high-frequency factor model, only the uniform convergence of the estimator for $(\mathbb{A}^{\dagger}\mathbb{A})^{-1/2}\mathbb{A}^{\dagger}$ can be derived.

The introduction of the magnitude matrices provides an insight into how larger noise-to-signal ratio can worsen the estimation. We can see that the convergence rates are affected by the magnitudes of both \mathbf{D}_G and \mathbf{D}_V , by noticing that $\underline{\tau}_G^{*-}$ directly affects the convergence rate and that $\bar{\tau}_G^{*+}$ and $\bar{\tau}_V^{*+}$ affect a_{nd} . The convergence rates are fastest (i.e., $(\log d)^{1/2}n^{-1/2} + m_{w,nd}d^{-1/2}$) when $\bar{\tau}_G = \underline{\tau}_G = \bar{\tau}_V = -1/2$. When common components of the microstructure noise have a larger or smaller magnitude than that of the efficient prices, the estimators will be worse. Note that $n^{\bar{\tau}_G^{*+}-2\underline{\tau}_G^{*-}} \cdot a_{nd}$ (see Theorem 3.1 (iii)) grows with $\bar{\tau}_G$ when $\bar{\tau}_G > -1/2$, and thus it may not converge to zero. However, this is not a big problem, since the convergence of the estimator of common components still holds if we re-scale the data by a factor of order $n^{-\bar{\tau}_G^{*+}}$. Nevertheless, if we require $n^{\bar{\tau}_G^{*+}-2\underline{\tau}_G^{*-}} \cdot a_{nd} = o(1)$, we only need to extend Assumption 8(i) to $n^{1+4\underline{\tau}_G^{*-}-2\bar{\tau}_G^{*+}-2\bar{\tau}_V^{*+}-2(\bar{\tau}_G^{*+}\vee\bar{\tau}_V^{*+})/(\log d) \to \infty$.

3.2. Second step: PCA estimation in cumulative form

The second step aims to obtain a $K_H \times K_G$ matrix $\hat{\boldsymbol{\beta}}$ such that $\hat{\boldsymbol{k}}^* \hat{\boldsymbol{\beta}}$ estimates the factors for microstructure noise. Denote $\boldsymbol{\beta} = (\mathbf{O}_{K_G \times K_F} \ \mathbf{I}_{K_G})^{\mathsf{T}}$ and $\boldsymbol{\beta}_{\perp} = (\mathbf{I}_{K_F} \ \mathbf{O}_{K_F \times K_G})^{\mathsf{T}}$. Then $\boldsymbol{g} := \boldsymbol{\hbar} \boldsymbol{\beta}$ is the matrix of true factors for the difference of microstructure noise, and $\boldsymbol{\ell} := \boldsymbol{\hbar} \boldsymbol{\beta}_{\perp}$ is the matrix of true factors for the difference of efficient prices. However, due to rotational indeterminacy of the estimated factors, $\hat{\boldsymbol{k}}^*$, instead of applying $\boldsymbol{\beta}$ and $\boldsymbol{\beta}_{\perp}$ directly on $\hat{\boldsymbol{k}}^*$ to obtain estimates of \boldsymbol{g} and $\boldsymbol{\ell}$, we need to find proper rotations of $\boldsymbol{\beta}$ and $\boldsymbol{\beta}_{\perp}$ that achieve this goal.

There are different ways to estimate such rotations of β and β_{\perp} . For example, Barigozzi et al. (2020b) use Johansen (1995)'s reduced rank estimation in a dynamic factor model to estimate the cointegration coefficients of non-stationary factors (also see Section 4.2). However, their method requires the specification of a finite-order vector autoregression (or a vector error correction model) prior to estimation. By contrast, we will use a second-step PCA on cumulated factors to estimate β and β_{\perp} , which allows for more general (non)stationarity structure of the underlying processes. This second step is similar to PANIC due to Bai and Ng (2004), which is developed from Stock and Watson (1988) and Harris (1997).

For
$$1 \leq s \leq n$$
, let $\widehat{H}_{s\Delta}^* = \sum_{s_1=1}^s \widehat{h}_{s_1\Delta}^*$ be an estimator of $H_{s\Delta}$. Define the demeaned $\widehat{H}_{s\Delta}^*$
as $\widehat{H}_{s\Delta}^{*c} = \widehat{H}_{s\Delta}^* - \overline{\widehat{H}^*}$, where $\overline{\widehat{H}^*} = n^{-1} \sum_{s=1}^n \widehat{H}_{s\Delta}^*$. In the matrix form, this can be written as
 $\widehat{\mathscr{H}^{*c}} = \widehat{\mathscr{H}^*} - \overline{\widehat{\mathscr{H}^*}}$, with $\widehat{\mathscr{H}^{*c}} = (\widehat{H}_{\Delta}^{*c}, \dots, \widehat{H}_{n\Delta}^{*c})^{\mathsf{T}}$, $\widehat{\mathscr{H}^*} = (\widehat{H}_{\Delta}^*, \dots, \widehat{H}_{n\Delta}^*)^{\mathsf{T}}$, and $\overline{\widehat{\mathscr{H}^*}} = \mathbf{1}_n \overline{\widehat{H}^*}^{\mathsf{T}}$. Define
the $K_H \times K_H$ matrix $\widehat{W} = n^{-1} \widehat{\mathscr{H}^{*c}} \widehat{\mathscr{H}^{*c}}$.

Let $\widehat{\boldsymbol{\beta}}_{\perp}$ be the matrix of eigenvectors associated with the largest K_F eigenvalues of \widehat{W} and let $\widehat{\boldsymbol{\beta}}$ be the matrix of eigenvectors associated with the rest of the K_G eigenvalues. Then, $\widehat{\boldsymbol{\beta}}_{\perp}^{\mathsf{T}} \widehat{\boldsymbol{h}}_t^*$ is an estimator of $\boldsymbol{\beta}_{\perp}^{\mathsf{T}} \boldsymbol{h}_t = \boldsymbol{f}_t$, and $\widehat{\boldsymbol{\beta}}^{\mathsf{T}} \widehat{\boldsymbol{h}}_t^*$ is an estimator of $\boldsymbol{\beta}_{\perp}^{\mathsf{T}} \boldsymbol{h}_t = \boldsymbol{f}_t$.

Lemma 3.1. Under the assumptions of Theorem 3.1, $\hat{\beta}_{\perp}$ and $\hat{\beta}$ are super-consistent in the sense that:

$$\widehat{\boldsymbol{\beta}} - \boldsymbol{\Xi}^{\mathsf{T}} \boldsymbol{\beta} \mathbf{Q}_{\boldsymbol{\beta}} = O_P(n^{-1}) \tag{3.9}$$

$$\widehat{\boldsymbol{\beta}}_{\perp} - \boldsymbol{\Xi}^{-1} \boldsymbol{\beta}_{\perp} \mathbf{Q}_{\boldsymbol{\beta}_{\perp}} = O_P(n^{-1}), \qquad (3.10)$$

where

$$\begin{aligned} \mathbf{Q}_{\boldsymbol{\beta}} &= [\boldsymbol{\beta}^{\mathsf{T}} \boldsymbol{\Xi} \boldsymbol{\Xi}^{\mathsf{T}} \boldsymbol{\beta}]^{-1} \boldsymbol{\beta}^{\mathsf{T}} \boldsymbol{\Xi} \widehat{\boldsymbol{\beta}}, \qquad \mathbf{Q}_{\boldsymbol{\beta}_{\perp}} = [\boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\boldsymbol{\Xi}^{\mathsf{T}})^{-1} \boldsymbol{\Xi}^{-1} \boldsymbol{\beta}_{\perp}]^{-1} \boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\boldsymbol{\Xi}^{\mathsf{T}})^{-1} \widehat{\boldsymbol{\beta}}_{\perp}, \\ \boldsymbol{\Xi} &= d^{1/2} \boldsymbol{\Lambda}_{H}^{\mathsf{T}} \widehat{\boldsymbol{\Lambda}}_{H} \widehat{\mathbf{D}}_{x, K_{H}}^{-1/2} = (\boldsymbol{\mathbb{A}}^{\mathsf{T}} \boldsymbol{\mathbb{A}})^{1/2} \mathbf{R}^{*}. \end{aligned}$$

The lemma shows that $\hat{\boldsymbol{\beta}}$ estimates a basis for the space spanned by $\boldsymbol{\Xi}^{\mathsf{T}}\boldsymbol{\beta}$. Using Lemma 3.1 and Theorem 3.1, we can prove the following theorem, which gives the convergence rate for the estimators of the factors for microstructure noise and the factor loadings for efficient prices. To this end, we define

$$\widehat{\mathbf{f}}^* = \widehat{\mathbf{\hbar}}^* \widehat{\mathbf{eta}}_{\perp}, \quad \widehat{\mathbf{g}}^* = \widehat{\mathbf{\hbar}}^* \widehat{\mathbf{eta}}, \quad \widehat{\mathbf{\Lambda}}^*_F = \widehat{\mathbf{\Lambda}}^*_H \widehat{\mathbf{eta}}_{\perp}, \quad ext{and} \quad \widehat{\mathbf{\Lambda}}^*_G = \widehat{\mathbf{\Lambda}}^*_H \widehat{\mathbf{eta}}_H$$

Theorem 3.2. Suppose that Assumptions 1-8 are satisfied. We have (i)

$$\left\|\widehat{\boldsymbol{f}}^* - \boldsymbol{f}\boldsymbol{\beta}_{\perp}^{\mathsf{T}}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\widehat{\boldsymbol{\beta}}_{\perp}\right\|_{\mathsf{MAX}} \leq \left\|\widehat{\boldsymbol{f}}^* - \boldsymbol{f}\boldsymbol{\beta}_{\perp}^{\mathsf{T}}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\widehat{\boldsymbol{\beta}}_{\perp}\right\|_{\mathsf{sp}} = O_P\left(n^{-2\underline{\tau}_G^{*-}} \cdot a_{nd}\right); \tag{3.11}$$

(ii)

$$\|\widehat{\boldsymbol{g}}^* - \boldsymbol{g}\mathbf{Q}_{\boldsymbol{\beta}}\|_{\mathsf{MAX}} \le \|\widehat{\boldsymbol{g}}^* - \boldsymbol{g}\mathbf{Q}_{\boldsymbol{\beta}}\|_{\mathsf{sp}} = O_P\left(n^{-2\underline{\tau}_G^{*-}} \cdot a_{nd}\right); \qquad (3.12)$$

(iii)

$$\left\|\widehat{\mathbf{\Lambda}}_{F}^{*}-\mathbf{\Lambda}_{F}\mathbf{Q}_{\boldsymbol{\beta}_{\perp}}\right\|_{\mathsf{MAX}}=O_{P}\left(n^{-2\underline{\tau}_{G}^{*-}}\cdot a_{nd}\right);$$
(3.13)

$$\left\|\widehat{\mathbf{\Lambda}}_{G}^{*}-\mathbf{\Lambda}_{G}\mathbf{D}_{G}\boldsymbol{\beta}^{\mathsf{T}}\Xi\widehat{\boldsymbol{\beta}}\right\|_{\mathsf{MAX}}=O_{P}\left(n^{-2\underline{\tau}_{G}^{*-}}\cdot a_{nd}\right),\tag{3.14}$$

where $\mathbf{Q}_{\boldsymbol{\beta}}$, $\mathbf{Q}_{\boldsymbol{\beta}_{\perp}}$ and $\boldsymbol{\Xi}$ are defined in Lemma 3.1.

This lemma shows that $\hat{\ell}^*$ and \hat{g}^* are estimators of the factors for the first-differenced efficient prices and microstructure noise with rotations $\beta_{\perp}^{\mathsf{T}}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\hat{\boldsymbol{\beta}}_{\perp}$ and $\mathbf{Q}_{\boldsymbol{\beta}}$, respectively, and that $\hat{\boldsymbol{\Lambda}}_F^*$ and $\hat{\boldsymbol{\Lambda}}_G^*$ are estimators of the factor loadings for efficient prices and microstructure noise with rotations $\mathbf{Q}_{\boldsymbol{\beta}_{\perp}}$ and $\boldsymbol{\beta}^{\mathsf{T}} \boldsymbol{\Xi} \hat{\boldsymbol{\beta}}$, respectively. When we estimate the first-differenced common components for microstructure noise and efficient prices (i.e., $\boldsymbol{g} \boldsymbol{\Lambda}_G^{\mathsf{T}}$ and $\boldsymbol{\ell} \boldsymbol{\Lambda}_F^{\mathsf{T}}$), these rotations cancel out. Thus, we have the following corollary.

Corollary 3.1. Suppose that Assumptions 1–8 are satisfied. We have (i)

$$\left|\widehat{\boldsymbol{f}}^{*}\widehat{\boldsymbol{\Lambda}}_{F}^{*^{\mathsf{T}}} - \boldsymbol{f}\boldsymbol{\Lambda}_{F}^{\mathsf{T}}\right\|_{\mathsf{MAX}} = O_{P}\left(n^{-2\underline{\tau}_{G}^{*^{-}}} \cdot a_{nd}\right);$$
(3.15)

(ii)

$$\left\|\widehat{g}^{*}\widehat{\Lambda}_{G}^{*^{\mathsf{T}}} - g\mathbf{D}_{G}\Lambda_{G}^{\mathsf{T}}\right\|_{\mathsf{MAX}} = O_{P}\left(n^{\overline{\tau}_{G}^{*} - 2\underline{\tau}_{G}^{*^{-}}} \cdot a_{nd}\right),\tag{3.16}$$

where $\bar{\tau}_{G}^{*} = 1/2 + \bar{\tau}_{G}$ and $\underline{\tau}_{G}^{*-} = (1/2 + \underline{\tau}_{G})_{-}$.

3.3. Third step: Cumulation of factors

Now we construct the estimates of \mathscr{F} and \mathscr{G} by cumulating $\widehat{\mathscr{F}}^* = \left(\widehat{f}^*_{\Delta}, \dots, \widehat{f}^*_{n\Delta}\right)^{\mathsf{r}}$ and $\widehat{\mathscr{Q}}^* = \left(\widehat{g}^*_{\Delta}, \dots, \widehat{g}^*_{n\Delta}\right)^{\mathsf{r}}$, respectively. That is, $\widehat{\mathscr{F}}^* = \left(\widehat{F}^*_{\Delta}, \dots, \widehat{F}^*_{n\Delta}\right)^{\mathsf{r}}$, $\widehat{\mathscr{G}}^* = \left(\widehat{G}^*_{\Delta}, \dots, \widehat{G}^*_{n\Delta}\right)^{\mathsf{r}}$, and $\widehat{\mathscr{H}}^* = \left(\widehat{H}^*_{\Delta}, \dots, \widehat{H}^*_{n\Delta}\right)^{\mathsf{r}}$, where

$$\widehat{F}_t^* = \sum_{s:0 \le s\Delta \le t} \widehat{f}_{s\Delta}^* \quad \text{and} \quad \widehat{G}_t^* = \sum_{s:0 \le s\Delta \le t} \widehat{g}_{s\Delta}^*, \quad \widehat{H}_t^* = \sum_{s:0 \le s\Delta \le t} \widehat{h}_{s\Delta}^*$$

More compactly, we can write $\mathscr{F} = \mathbf{L}_n \mathbf{f}$, $\mathscr{G} = \mathbf{L}_n \mathbf{g}$ and $\mathscr{H} = \mathbf{L}_n \mathbf{h}$, and their estimates, $\widehat{\mathscr{F}}^* = \mathbf{L}_n \widehat{\mathbf{f}}^*$, $\widehat{\mathscr{G}}^* = \mathbf{L}_n \widehat{\mathscr{G}}^*$ and $\widehat{\mathscr{H}}^* = \mathbf{L}_n \widehat{\mathscr{H}}^*$ for some matrix \mathbf{L}_n consisting of zeros and ± 1 . Since $\|\mathbf{L}_n\|_{sp} = 1$, the convergence rates for the estimators of the first-differenced factors are also the convergence rates for the aggregated factors. This leads to the following theorem.

Theorem 3.3. Suppose that Assumptions 1–8 are satisfied. Then, (i) for factors, we have

$$\left\|\widehat{\mathscr{H}^*} - (\mathscr{H} - \mathbf{1}_n \boldsymbol{H}_0^{\mathsf{T}})(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\right\|_{\mathsf{MAX}} \le \left\|\widehat{\mathscr{H}^*} - (\mathscr{H} - \mathbf{1}_n \boldsymbol{h}_0^{\mathsf{T}})(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\right\|_{\mathsf{sp}} = O_P\left(n^{-2\underline{\tau}_G^{*-}} \cdot a_{nd}\right), \quad (3.17)$$

$$\left\|\widehat{\mathscr{G}}^* - (\mathscr{G} - \mathbf{1}_n \mathbf{G}_0^{\mathsf{T}}) \mathbf{Q}_{\boldsymbol{\beta}}\right\|_{\mathsf{MAX}} \leq \left\|\widehat{\mathscr{G}}^* - (\mathscr{G} - \mathbf{1}_n \mathbf{G}_0^{\mathsf{T}}) \mathbf{Q}_{\boldsymbol{\beta}}\right\|_{\mathsf{sp}} = O_P\left(n^{-2\underline{\tau}_G^{*-}} \cdot a_{nd}\right),\tag{3.18}$$

and

$$\left\|\widehat{\mathscr{F}}^{*} - (\mathscr{F} - \mathbf{1}_{n} \mathbf{F}_{0}^{\mathsf{T}}) \boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\boldsymbol{\Xi}^{\mathsf{T}})^{-1} \widehat{\boldsymbol{\beta}}_{\perp} \right\|_{\mathsf{MAX}} \leq \left\|\widehat{\mathscr{F}}^{*} - (\mathscr{F} - \mathbf{1}_{n} \mathbf{F}_{0}^{\mathsf{T}}) \boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\boldsymbol{\Xi}^{\mathsf{T}})^{-1} \widehat{\boldsymbol{\beta}}_{\perp} \right\|_{\mathsf{sp}} = O_{P} \left(n^{-2\underline{\tau}_{G}^{*-}} \cdot a_{nd} \right);$$
(3.19)

(ii) for common components, we have

$$\left\|\widehat{\mathscr{H}}^{*}(\widehat{\boldsymbol{\Lambda}}_{H}^{*})^{\mathsf{T}} - (\mathscr{H} - \mathbf{1}_{n}\boldsymbol{H}_{0}^{\mathsf{T}})\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\right\|_{\mathsf{MAX}} = O_{P}\left(n^{\overline{\tau}_{G}^{*+} - 2\underline{\tau}_{G}^{*-}} \cdot a_{nd}\right),\tag{3.20}$$

$$\left\|\widehat{\mathscr{G}}^{*}(\widehat{\boldsymbol{\Lambda}}_{G}^{*})^{\mathsf{T}} - (\mathscr{G} - \mathbf{1}_{n}\boldsymbol{G}_{0}^{\mathsf{T}})\mathbf{D}_{G}\boldsymbol{\Lambda}_{G}^{\mathsf{T}}\right\|_{\mathsf{MAX}} = O_{P}\left(n^{\overline{\tau}_{G}^{*} - 2\underline{\tau}_{G}^{*-}} \cdot a_{nd}\right),\tag{3.21}$$

and

$$\left\|\widehat{\mathscr{F}}^{*}(\widehat{\mathbf{\Lambda}}_{F}^{*})^{\mathsf{T}} - (\mathscr{F} - \mathbf{1}_{n}\mathbf{F}_{0}^{\mathsf{T}})\mathbf{\Lambda}_{F}^{\mathsf{T}}\right\|_{\mathsf{MAX}} = O_{P}\left(n^{-2\underline{\tau}_{G}^{*-}} \cdot a_{nd}\right),\tag{3.22}$$

where \mathbf{Q}_{β} , $\mathbf{Q}_{\beta_{\perp}}$ and Ξ are defined in Lemma 3.1.

When $\bar{\tau}_G = \underline{\tau}_G = \bar{\tau}_V = -1/2$ and $m_{w,nd} = O(1)$, the uniform convergence rate of $\widehat{\mathscr{H}}^*$ is $O_P((\log d)^{1/2}n^{-1/2} + d^{-1/2})$. In comparison, the uniform consistency result for the low-frequency factor model in Bai and Ng (2004) is $O_P(n^{-3/4} + d^{-1/2})$ (see Lemma 2 in Bai and Ng (2004)). Our estimators have slower convergence rates, which is mainly due to the different settings.

4. Simulation

4.1. Number of factors

In this paper, we apply the commonly-used information criterion proposed by Bai and Ng (2002) to estimate the total number of factors. Then, we use the PANIC test procedure proposed by Bai and Ng (2004) to determine the number of factors for efficient prices, and compare it to Hallin and Liška (2007)'s spectral method.

To introduce Bai and Ng (2002)'s information criterion, we denote by \overline{Q} a finite positive integer that is no smaller than K_H . For any $1 \leq q_H \leq \overline{Q}$, we let $\hat{\mathcal{R}}^*(q_H) = \left(\hat{\mathbf{h}}^*_{\Delta}(q_H), \ldots, \hat{\mathbf{h}}^*_{n\Delta}(q_H)\right)^{\mathsf{T}}$ be the matrix of estimated factors when the total number of factors is assumed to be q_H , and denote by $\hat{\mathbf{A}}^*_H(q_H)$ the corresponding loadings matrix. The information criterion is defined as

$$\mathsf{IC}_1(q_H) = \log\left[\mathsf{V}_n(q_H)\right] + q_H\left(\frac{n+d}{nd}\right)\log(\frac{nd}{n+d}),\tag{4.1}$$

where $V_n(q_H) = \|x - \hat{k}^*(q_H)(\widehat{\Lambda}_H^*(q_H))^{\mathsf{T}}\|_F$. The total number of factors is then estimated as

$$\widehat{K}_H = \underset{0 \le q_H \le \overline{Q}}{\arg\min} \mathsf{IC}_1(q_H), \tag{4.2}$$

with $\mathsf{IC}_1(0) = ||x||_F$ for convention. For consistency of \widehat{K}_H , we refer to the asymptotic results given in Theorem 2 of Bai and Ng (2002).

Bai and Ng (2004) propose two tests to determine the number of nonstationary factors. We adopt the one that does not specify a finite order VAR representation. For any $1 \leq q_F \leq \hat{K}_H$, we let $\widehat{\mathscr{F}}^*(q_F) = \widehat{\mathscr{H}}^* \widehat{\beta}_{\perp}(q_F)$, when the number of factors for efficient prices is assumed to be q_F . Let $\widehat{\xi}_t^F$ be the residuals from estimating a first-order VAR of $\widehat{\mathscr{F}}^*(q_F)$ and $\widehat{\Sigma}_{L,F}(q_F)$ be the estimated long-run covariance matrix of $\widehat{\xi}_t^F$. The test statistic for $H_0: K_F = q_F$ is defined by

$$MQ(q_F) = n(\nu(q_F) - 1)$$
 (4.3)

where $\nu(q_F)$ is the smallest eigenvalue of

$$\left[\sum_{s=2}^{n} \frac{1}{2} \left(\widehat{F}_{s\Delta}^{*}(q_F) \widehat{F}_{(s-1)\Delta}^{*\dagger}(q_F) + \widehat{F}_{(s-1)\Delta}^{*}(q_F) \widehat{F}_{s\Delta}^{*\dagger}(q_F) \right) - n \widehat{\Sigma}_{L,F}(q_F) \right] \left(\sum_{s=2}^{n} \widehat{F}_{s\Delta}^{*}(q_F) \widehat{F}_{s\Delta}^{*\dagger}(q_F) \right)^{-1}.$$
(4.4)

For a given significance level α , define

$$\widehat{K}_F = \max_{0 \le q_F \le \widehat{K}_H, \ p_{MQ}(q_F) > \alpha} q_F, \tag{4.5}$$

where $p_{MQ}(q_F)$ is the p-value of the statistic $MQ(q_F)$ obtained by simulation using vector standard Brownian motions, and we define $p_{MQ}(0) = 1$ for convention. We refer the reader to Theorem 1 of Bai and Ng (2004) for the asymptotic distribution of the test statistic.

Hallin and Liška (2007)'s spectral method can also be used to estimate the number of factors for efficient prices. Let $\hat{\Gamma}_k$ be the $d \times d$ sample lag-k autocovariance matrix of \boldsymbol{x}_t . Define the lag window estimator of the spectral density matrix of \boldsymbol{x}_t by

$$\widehat{\Sigma}_x(\theta) = \frac{1}{2\pi} \sum_{k=-B_n}^{B_n} \widehat{\Gamma}_k e^{-ik\theta} w(B_n^{-1}k), \qquad (4.6)$$

where B_n is a suitable bandwidth and $w(\cdot)$ is a positive even weight function. Let $\nu_l(\theta)$ be the *l*-th largest eigenvalue of $\widehat{\Sigma}_x(\theta)$ and define the following information criteria

$$\mathsf{IC}_{2,i}(q) = \log\left(\frac{1}{n(2B_n+1)}\sum_{h=-B_n}^{B_n}\sum_{l=q+1}^d\nu_l(\theta_h)\right) + qs_i(n,d), \quad i = 1, 2, 3,$$
(4.7)

and

$$\mathsf{IC}_{3,i}(q) = \log\left(\frac{1}{n}\sum_{l=q+1}^{d}\nu_l(0)\right) + qs_i(n,d), \quad i = 1, 2, 3,$$
(4.8)

where, with $M_n = \lfloor 0.75\sqrt{n} \rfloor$ and a constant ι ,

$$s_1(n,d) = \iota \cdot (M_n^{-2} + M_n^{1/2}n^{-1/2} + d^{-1}) \cdot \log(\min\{M_n^2, M_n^{-1/2}n^{1/2}, d\})$$
$$s_2(n,d) = \iota \cdot (\min\{M_n^2, M_n^{-1/2}n^{1/2}, d\})^{-1/2},$$

and

$$s_3(n,d) = \iota \cdot (\min\{M_n^2, M_n^{-1/2}n^{1/2}, d\})^{-1} \cdot \log(\min\{M_n^2, M_n^{-1/2}n^{1/2}, d\}).$$

Barigozzi et al. (2020b) show that K_H can be estimated as

$$\widehat{K}_{H} = \underset{1 \le q \le \overline{Q}}{\arg\min} \mathsf{IC}_{2.i}(q) \tag{4.9}$$

for any i = 1, 2, 3, and K_F can be estimated as

$$\widehat{K}_F = \arg\min_{1 \le q \le \overline{Q}} \mathsf{IC}_{3.i}(q) \tag{4.10}$$

for any i = 1, 2, 3. In practice, one can let $\overline{Q} = \widehat{K}_H$ in (4.10) to make sure $\widehat{K}_F \leq \widehat{K}_H$.

4.2. Alternative approaches for comparison

We consider two alternative approaches for comparison. The first approach, denoted as DPCA, estimates the factors as $\hat{\mathcal{R}}$ in (3.2) instead of $\hat{\mathcal{R}}^*$, and uses $\hat{\mathcal{R}}$ for the second-step PCA while keeping the rest of the steps exactly the same. The second approach is proposed by Barigozzi et al. (2020a) and Barigozzi et al. (2020b) and is denoted as PCA*-VECM. The PCA*-VECM uses \hat{h}_t^* from the first-step PCA* to construct a Vector Error Correction Model (VECM) in order to estimate β and β_{\perp} , while in our method, we use a second-step PCA to estimate β and β_{\perp} . We set the lag of VECM to 1 for simplicity, and for a lag length larger than 1, we refer the reader to Chapter 6 of Johansen (1995). Specifically, the PCA*-VECM uses Johansen (1995)'s reduced rank regression method to estimate a VECM for \hat{h}_t^* :

Step 1: Implement OLS of \hat{h}_t^* and $\hat{H}_{t-\Delta}^*$ on $\hat{h}_{t-\Delta}^*$ to get residuals $\hat{e}_{0,t}$ and $\hat{e}_{1,t}$, respectively.

Step 2: Let $\widehat{\mathbf{S}}_{ij} = n^{-1} \sum_{s=1}^{n} \widehat{e}_{i,s\Delta} \widehat{e}_{j,s\Delta}^{\mathsf{T}}$ for i, j = 0, 1. Then let $\widehat{\boldsymbol{\beta}} = (\widehat{\boldsymbol{\beta}}_{1}, \dots, \widehat{\boldsymbol{\beta}}_{K_{G}})$, where $\widehat{\boldsymbol{\beta}}_{l}, l = 1, \dots, K_{G}$, is the eigenvector belonging to the *l*-th largest eigenvalue of the matrix $(\widehat{\mathbf{S}}_{11} - \widehat{\mathbf{S}}_{10}\widehat{\mathbf{S}}_{00}^{-1}\widehat{\mathbf{S}}_{01})$. Define $\widehat{\boldsymbol{\beta}}_{\perp}$ as the orthogonal complement matrix of $\widehat{\boldsymbol{\beta}}$ such that $\widehat{\boldsymbol{\beta}}_{\perp}^{\mathsf{T}}\widehat{\boldsymbol{\beta}} = \mathbf{O}_{K_{F} \times K_{G}}$ and $\widehat{\boldsymbol{\beta}}_{\perp}^{\mathsf{T}}\widehat{\boldsymbol{\beta}}_{\perp} = \mathbf{I}_{K_{F}}$. Note that even if the estimated factors are not normalised in Step 1, the residuals $\widehat{\boldsymbol{e}}_{0,t}$ and $\widehat{\boldsymbol{e}}_{1,t}$ will

Note that even if the estimated factors are not normalised in Step 1, the residuals $e_{0,t}$ and $e_{1,t}$ will not change and therefore the estimates of $\hat{\beta}_{\perp}$ and $\hat{\beta}$ will not be affected. Also note that the factors for efficient prices follow a diffusion model and hence, are heteroskedastic. One can implement more efficient estimation of VECM under heteroskedasticity (e.g., generalized least squares estimation in Seo (2007) and Herwartz and Lütkepohl (2011)). But we do not pursue this in our paper.

4.3. Data generating processes

The data generating process in our simulation is similar to that in Bollerslev et al. (2019). The observable prices are the sum of efficient prices and microstructure noise. The former has two orthogonal factors and the latter has one factor.

The two factors in efficient prices are independently generated from a GARCH diffusion model as in Andersen and Bollerslev (1998),

$$\begin{cases} f_{it} = \sigma_{fit} \mathsf{d}B_{it}, \\ \mathsf{d}\sigma_{fit}^2 = \kappa_{fi}(\theta_{fi} - \sigma_{fit}^2) \mathsf{d}t + \lambda_{fi}\sigma_{fit}^2 \mathsf{d}W_{it}, \end{cases}$$
(4.11)

for i = 1, 2, where B_{it} and W_{it} are dependent Brownian motions with $\operatorname{corr}(B_{it}, W_{it}) = -0.5$. The parameters are set as $\kappa_{f1} = \kappa_{f2} = 0.035$, $\theta_{f1} = 0.636$, $\theta_{f2} = 0.3$, $\lambda_{fi} = \sqrt{2\kappa_{fi}\phi_{fi}}$, $\phi_{f1} = \phi_{f2} = 0.296$, and initial value $(f_{i0}, \sigma_{fi0}^2) = (0, \theta_{fi})$. Then we draw the factor loadings of efficient prices independently from a normal distribution with mean zero and unit variance.

The idiosyncratic components of efficient prices are generated as $U_{it} = \sigma_{it} \mathsf{W}_{it}^U$, where W_{it}^U is a Brownian motion and σ_{it} is generated by three different models for different *i*.

• For $1 \le i \le \lfloor d/3 \rfloor$, the volatility process is generated by an exponential ARCH diffusion limit model as in Nelson (1990):

$$\mathsf{d}\log(\sigma_{it}^2) = -0.6(0.157 - \log(\sigma_{it}^2))\mathsf{d}t + 0.25\mathsf{d}B_{it}^U \tag{4.12}$$

with initial value $\log(\sigma_{i0}^2) = 0.157$, where B_{it}^U is a Brownian motion with $\operatorname{corr}(B_{it}^U, W_{it}^U) = -0.3$.

• For $\lfloor d/3 \rfloor + 1 \le i \le \lfloor 2d/3 \rfloor$, the volatility process is generated by a GARCH-M diffusion limit model as in Nelson (1990),

$$\mathsf{d}(\sigma_{it}^2) = (0.1 - \sigma_{it}^2)\mathsf{d}t + 0.2\sigma_{it}^2\mathsf{d}B_{it}^U$$
(4.13)

with initial value $\sigma_{i0}^2 = 0.1$, where B_{it}^U is a Brownian motion with $\operatorname{corr}(B_{it}^U, W_{it}^U) = -0.3$.

• For $\lfloor 2d/3 \rfloor + 1 \leq i \leq d$, the volatility process is generated by a GARCH diffusion model as in Andersen and Bollerslev (1998),

$$\mathsf{d}(\sigma_{it}^2) = 0.035(0.636 - \sigma_{it}^2)\mathsf{d}t + 0.2\sigma_{it}^2\mathsf{d}B_{it}^U, \tag{4.14}$$

with initial value $\sigma_{i0}^2 = 0.636$, where B_{it}^U is a Brownian motion with $\operatorname{corr}(B_{it}^U, W_{it}^U) = -0.3$.

The two dimensional Brownian motion (B_{it}^U, W_{it}^U) is independent over $1 \le i \le d$ and also independent with the driving Brownian motions (B_{1t}, W_{1t}) and (B_{2t}, W_{2t}) for factors of efficient prices.

As for microstructure noise, we introduce the noise-to-signal ξ_G^2 and ξ_V^2 as in Bollerslev et al. (2019), which take values $n^{2\bar{\tau}_G}$ and $n^{2\bar{\tau}_V}$, respectively, with $\bar{\tau}_G = -0.4$, -0.6 and $\bar{\tau}_V = -0.4$, -0.6. The variance of the factor for microstructure noise satisfies $\operatorname{Var}(G_t) = 0.5\xi_G^2(\frac{1}{nd}\sum_{i=1}^d \sum_{t=1}^n \sigma_{*,it}^4)^{1/2}$, and is thus time-invariant, where $\sigma_{*,it}$ is the spot volatility of the efficient price process of asset *i* at time *t*. The variance of idiosyncratic component V_{it} makes up $0.1\xi_V^2$ of the total variance, that is $\operatorname{Var}(V_{it}) = 0.5\xi_G^2(\frac{1}{nd}\sum_{i=1}^d \sum_{t=1}^n \sigma_{*,it}^4)^{1/2}$.

 $0.1\xi_V^2(\frac{1}{n}\sum_{t=1}^n \sigma_{*,it}^4)^{1/2}$. We draw the factor G_t independently from a normal distribution with mean zero and variance $\mathsf{Var}(G_t)$, and draw V_{it} independently from a normal distribution with mean zero and variance $\mathsf{Var}(V_{it})$. Finally, we draw the factor loadings, λ_{Gi} , i = 1, ..., d, of microstructure noise independently from a normal distribution with mean one and unit variance.

We assume that the prices are synchronously recorded once every one or five minutes during 6.5 trading hours, that is n = 390 or 78. The number of assets is assumed to be d = 50, 100, 300 and 500. We present the simulation results based on 1000 Monte Carlo replications.

4.4. Simulation results

Firstly, we provide simulation results for the estimation of number of factors by the information criteria described in Section 4.1, as well as by the PANIC test with different significance levels, 1%, 5% and 10%. More specifically, we use Bai and Ng (2002)'s information criterion, IC₁, to estimate the total number of factors and then use the PANIC test to identify the number of factors for efficient prices. We compare this against Hallin and Liška (2007)'s information criteria IC_{2,i} and IC_{3,i}, i = 1, 2, 3. For the PANIC test in (4.5), \hat{K}_H is determined by IC₁. We set $\overline{Q} = 10$ in both (4.9) and (4.10), and in Hallin and Liška (2007)'s information criteria, we set $B_n = 100$, $\iota = 0.5$, and w(x) = 1 - |x| (i.e., the Bartlett weight function).

Tables 2 and 3 provide the average number of factors determined by each method (over 1000 replications) for n = 78 and n = 390, respectively. It can be seen that IC_1 has excellent performance in estimating the total number of factors in all scenarios. Other information criteria perform differently. Notice that in Section 4.1, $IC_{2,i}$, i = 1, 2, 3, are used to estimate the total number of factors rather than factors for efficient prices. When $\bar{\tau}_V = -0.6$, i.e., when the idiosyncratic error of microstructure noise is relatively small, the average estimated number of factors by $IC_{2,i}$, i = 1, 2, 3, is close to 2. When $\bar{\tau}_V$ increases to -0.4, this number decreases, especially for $\mathsf{IC}_{2.3}$ when $\bar{\tau}_G = -0.6$ (i.e., when the noise-to-signal ratio is relatively high). Indeed when $\bar{\tau}_G = -0.6$ and $\bar{\tau}_V = -0.4$, the average number of factors estimated from $|C_{2,3}|$ falls well below 1. The same pattern applies to $|C_{3,i}|$, i = 1, 2, 3. The results show that when used to estimate the total number of factors, the criteria $IC_{2,i}$, i = 1, 2, 3, underestimate it in all scenarios. But similar to $IC_{3,i}$, i = 1, 2, 3, they provide a good estimate of the number of factors for efficient prices when the idiosyncratic error of microstructure noise is relatively small. However, when the magnitude of the idiosyncratic error for microstructure noise increases, $IC_{2.3}$ and $IC_{3.3}$ seriously underestimate this number. The PANIC tests using the MQ statistic at different significance levels are denoted as $MQ_{1\%}$, $MQ_{5\%}$, and $MQ_{10\%}$ in Tables 2 and 3. They are used to determine the number of factors for efficient prices and perform satisfactorily in all scenarios, in particular for $MQ_{1\%}$. In summary, IC₁ is very satisfactory in determining the number of total factors, and the PANIC test with 1% significance level is the most robust method to determine the number of factors for efficient prices, outperforming $IC_{2.2}$, which is the best performer among all of Hallin and Liška (2007)'s information criteria.

Next, we compare the estimation of common components in the dual factor model. Our method is denoted as DPCA^{*}. For the alternative methods discussed in Section 4.2, we denote them as DPCA and PCA^{*}-VECM, respectively. We use the relative estimation error (REE) to measure the performance of different methods. It is defined as

$$\text{REE} = \|\mathbf{M} - \mathbf{M}\| / \|\mathbf{M}\|,$$

	IC_1	IC _{2.1}	$IC_{2.2}$	$IC_{2.3}$	IC _{3.1}	IC _{3.2}	IC _{3.3}	$MQ_{1\%}$	$MQ_{5\%}$	$MQ_{10\%}$
				$n = 78, \bar{\tau}$	G = -0.4,	$\bar{\tau}_V = -0.$	4			
d=50	3.026	1.910	1.976	1.419	1.874	1.984	1.477	1.998	1.931	1.855
	(0.025)	(0.082)	(0.023)	(0.260)	(0.110)	(0.064)	(0.284)	(0.030)	(0.086)	(0.150)
d=100	3.001	1.889	1.956	1.271	1.852	1.921	1.335	1.983	1.931	1.858
	(0.001)	(0.099)	(0.042)	(0.224)	(0.126)	(0.073)	(0.287)	(0.017)	(0.072)	(0.134)
d=300	3.000	1.889	1.951	1.183	1.846	1.924	1.195	1.986	1.942	1.871
	(0.000)	(0.099)	(0.047)	(0.176)	(0.130)	(0.070)	(0.239)	(0.014)	(0.059)	(0.122)
d = 500	3.000	1.888	1.945	1.147	1.853	1.932	1.176	1.988	1.940	1.869
	(0.000)	(0.100)	(0.052)	(0.162)	(0.126)	(0.063)	(0.227)	(0.012)	(0.062)	(0.130)
				$n = 78, \bar{\tau}$	$_{G} = -0.6,$	$\bar{\tau}_V = -0.$	4			
d=50	3.026	1.705	1.914	0.614	1.688	1.898	0.821	2.015	1.939	1.861
	(0.025)	(0.246)	(0.083)	(0.445)	(0.275)	(0.136)	(0.535)	(0.047)	(0.091)	(0.156)
d = 100	3.001	1.647	1.847	0.277	1.626	1.797	0.492	1.984	1.930	1.858
	(0.001)	(0.287)	(0.140)	(0.255)	(0.312)	(0.182)	(0.400)	(0.018)	(0.073)	(0.134)
d = 300	3.000	1.650	1.851	0.114	1.602	1.771	0.263	1.987	1.943	1.873
	(0.000)	(0.308)	(0.143)	(0.119)	(0.344)	(0.203)	(0.254)	(0.013)	(0.058)	(0.121)
d = 500	3.000	1.652	1.834	0.083	1.590	1.788	0.180	1.988	1.939	1.869
	(0.000)	(0.337)	(0.171)	(0.086)	(0.408)	(0.227)	(0.178)	(0.012)	(0.063)	(0.130)
		-		$n = 78, \bar{\tau}$	G = -0.4,	$\bar{\tau}_V = -0.$	6			
d=50	3.001	2.000	2.000	1.983	1.997	2.020	1.955	1.991	1.949	1.884
	(0.001)	(0.000)	(0.000)	(0.017)	(0.005)	(0.024)	(0.043)	(0.011)	(0.056)	(0.117)
d = 100	3.000	2.000	2.000	1.975	1.998	1.999	1.952	1.986	1.942	1.881
	(0.000)	(0.000)	(0.000)	(0.024)	(0.002)	(0.001)	(0.046)	(0.014)	(0.063)	(0.115)
d = 300	3.000	2.000	2.000	1.990	2.000	2.000	1.969	1.991	1.951	1.890
	(0.000)	(0.000)	(0.000)	(0.010)	(0.000)	(0.000)	(0.030)	(0.009)	(0.051)	(0.108)
d = 500	3.000	2.000	2.000	1.990	2.000	2.000	1.974	1.988	1.948	1.887
	(0.000)	(0.000)	(0.000)	(0.010)	(0.000)	(0.000)	(0.025)	(0.012)	(0.055)	(0.114)
				$n = 78, \bar{\tau}$	G = -0.6,	$\bar{\tau}_V = -0.$	6			
d=50	3.001	2.000	2.000	1.927	1.989	2.017	1.844	2.008	1.956	1.892
	(0.001)	(0.000)	(0.000)	(0.074)	(0.013)	(0.021)	(0.156)	(0.028)	(0.054)	(0.114)
d=100	3.000	1.999	2.000	1.915	1.992	1.996	1.847	1.987	1.945	1.881
	(0.000)	(0.001)	(0.000)	(0.086)	(0.008)	(0.004)	(0.152)	(0.015)	(0.060)	(0.115)
d = 300	3.000	2.000	2.000	1.951	1.998	1.999	1.909	1.992	1.951	1.890
	(0.000)	(0.000)	(0.000)	(0.059)	(0.002)	(0.001)	(0.099)	(0.008)	(0.051)	(0.108)
$\overline{d=500}$	3.000	2.000	2.000	1.951	2.000	2.000	1.914	1.988	1.949	1.887
	(0.000)	(0.000)	(0.000)	(0.057)	(0.000)	(0.000)	(0.109)	(0.012)	(0.054)	(0.114)

Table 2: Average estimated number of factors and standard deviation (in parentheses) for sampling frequency=5 mins, true number of factors=2 (efficient prices)+1 (microstructure noise)

	IC_1	IC _{2.1}	$IC_{2.2}$	$IC_{2.3}$	IC _{3.1}	$IC_{3.2}$	IC _{3.3}	$MQ_{1\%}$	$MQ_{5\%}$	$MQ_{10\%}$
				$n = 390, \bar{\tau}$	$\bar{r}_G = -0.4$	$, \bar{\tau}_V = -0$.4			
d = 50	3.035	1.907	1.976	1.508	1.860	1.947	1.541	1.989	1.901	1.815
	(0.036)	(0.084)	(0.023)	(0.250)	(0.123)	(0.050)	(0.269)	(0.027)	(0.101)	(0.177)
d=100	3.004	1.913	1.964	1.266	1.843	1.913	1.388	1.987	1.917	1.827
	(0.004)	(0.080)	(0.035)	(0.195)	(0.132)	(0.080)	(0.264)	(0.015)	(0.080)	(0.165)
d=300	3.000	1.928	1.954	1.133	1.846	1.894	1.246	1.987	1.920	1.850
	(0.000)	(0.067)	(0.044)	(0.115)	(0.130)	(0.095)	(0.198)	(0.013)	(0.080)	(0.144)
d = 500	3.000	1.924	1.967	1.076	1.846	1.896	1.176	1.988	1.923	1.850
	(0.000)	(0.070)	(0.032)	(0.070)	(0.130)	(0.093)	(0.155)	(0.012)	(0.071)	(0.146)
				$n = 390, \bar{\tau}$	$\bar{r}_G = -0.6$	$, \bar{\tau}_V = -0$.4			
d = 50	3.035	1.210	1.645	0.383	1.381	1.684	0.626	1.991	1.909	1.818
	(0.036)	(0.350)	(0.245)	(0.253)	(0.396)	(0.250)	(0.415)	(0.027)	(0.101)	(0.175)
d=100	3.004	1.034	1.395	0.081	1.201	1.453	0.317	1.990	1.920	1.827
	(0.004)	(0.281)	(0.291)	(0.075)	(0.419)	(0.358)	(0.253)	(0.010)	(0.078)	(0.165)
d=300	3.000	0.847	1.059	0.010	1.002	1.200	0.089	1.989	1.920	1.851
	(0.000)	(0.292)	(0.292)	(0.010)	(0.450)	(0.432)	(0.083)	(0.011)	(0.080)	(0.143)
d = 500	3.000	0.777	0.955	0.004	0.906	1.049	0.057	1.988	1.923	1.850
_	(0.000)	(0.306)	(0.319)	(0.004)	(0.440)	(0.443)	(0.054)	(0.012)	(0.071)	(0.146)
				$n = 390, \bar{\imath}$	$\bar{r}_G = -0.4$	$, \bar{\tau}_V = -0$.6			
d=50	3.000	2.000	2.000	2.000	1.999	1.999	1.991	1.993	1.931	1.847
	(0.000)	(0.000)	(0.000)	(0.000)	(0.001)	(0.001)	(0.009)	(0.007)	(0.072)	(0.154)
d=100	3.000	2.000	2.000	2.000	2.000	2.000	1.999	1.991	1.932	1.848
	(0.000)	(0.000)	(0.000)	(0.000)	(0.000)	(0.000)	(0.001)	(0.009)	(0.067)	(0.147)
d=300	3.000	2.000	2.000	2.000	2.000	2.000	2.000	1.989	1.926	1.863
	(0.000)	(0.000)	(0.000)	(0.000)	(0.000)	(0.000)	(0.000)	(0.011)	(0.073)	(0.132)
d = 500	3.000	2.000	2.000	2.000	2.000	2.000	2.000	1.990	1.931	1.860
	(0.000)	(0.000)	(0.000)	(0.000)	(0.000)	(0.000)	(0.000)	(0.010)	(0.064)	(0.137)
				$n = 390, \bar{\tau}$	$\bar{r}_G = -0.6$	$, \bar{\tau}_V = -0$.6			
d = 50	3.000	2.000	2.000	1.980	1.986	1.993	1.947	1.995	1.935	1.852
	(0.000)	(0.000)	(0.000)	(0.022)	(0.014)	(0.007)	(0.052)	(0.007)	(0.071)	(0.150)
d=100	3.000	2.000	2.000	1.996	1.998	1.999	1.961	1.991	1.935	1.849
	(0.000)	(0.000)	(0.000)	(0.004)	(0.002)	(0.001)	(0.038)	(0.009)	(0.065)	(0.146)
d=300	3.000	2.000	2.000	2.000	2.000	2.000	1.986	1.990	1.927	1.863
	(0.000)	(0.000)	(0.000)	(0.000)	(0.000)	(0.000)	(0.014)	(0.010)	(0.072)	(0.132)
d=500	3.000	2.000	2.000	2.000	2.000	2.000	1.995	1.989	1.931	1.861
	(0.000)	(0.000)	(0.000)	(0.000)	(0.000)	(0.000)	(0.005)	(0.011)	(0.064)	(0.136)

Table 3: Average estimated number of factors and standard deviation (in parentheses) for sampling frequency=1 min, true number of factors=2 (efficient prices)+1 (microstructure noise)

where $\|\cdot\|$ can be the Frobenius norm, the max norm or the spectral norm, and $\widehat{\mathbf{M}}$ is an estimate of the matrix \mathbf{M} , which varies from place to place (depending on the quantity being estimated). In Tables 4–7, "diff.total" refers to the total common component, $\Re \Lambda_H^{\mathsf{T}}$. "EP" and "diff.EP" refer to the common components for efficient prices and efficient returns. That is, $\mathrm{EP} = \mathscr{F} \Lambda_F^{\mathsf{T}}$ and diff.EP = $\mathscr{F} \Lambda_F^{\mathsf{T}}$, respectively. Similarly, "MN" and "diff.MN" refer to the common components for microstructure noise and its first difference, i.e., $\mathrm{MN} = \mathscr{G} \Lambda_G^{\mathsf{T}}$ and diff.MN = $\mathscr{G} \Lambda_G^{\mathsf{T}}$, respectively. In the following study, we will only consider the cases $\overline{\tau}_G = \overline{\tau}_V = -0.6$ and $\overline{\tau}_G = \overline{\tau}_V = -0.4$. To focus on the comparison of different methods in estimating the common components, we assume that the true numbers of factors for efficient prices and microstructure noise are known. Otherwise, one can obtain correct estimates of the numbers of factors by IC_1 and the PANIC test in most cases (as can be seen from Tables 2 and 3).

Table 4 gives the simulation results when n = 78, $\bar{\tau}_G = -0.6$, $\bar{\tau}_V = -0.6$. In this case, the microstructure noise is smaller than the efficient returns in magnitude. We can see that the REE in all columns decreases when d increases. The estimates of the total common component, $\hbar \Lambda_H^{\mathsf{T}}$, are the same for all the three methods, and the corresponding REE values are listed in the "diff.total" column. Since the estimates of $\not A_F^{\mathsf{T}}$ and $\mathcal{Q}\Lambda_G^{\mathsf{T}}$ are decomposed from the estimate of $\hbar \Lambda_H^{\mathsf{T}}$, the "diff.EP" and "diff.MN" have larger REEs than "diff.total" for all the three methods. Overall, DPCA* performs the best.

Table 5 presents the simulation results when n = 78, $\bar{\tau}_G = -0.4$, $\bar{\tau}_V = -0.4$. In this case, the microstructure noise is larger than the efficient returns in magnitude. Similar to Table 4, the REEs decrease when d increases. DPCA* performs the best, while DPCA performs the poorest. This may be due to the fact that DPCA* normalises the factors in the first step PCA, while DPCA does not. When the magnitude of the factors for microstructure noise, $\mathbf{D}_G \mathbf{g}_t$, is larger than that of the factors for efficient returns, \mathbf{f}_t , the magnitude of the cumulated factors for microstructure noise, $\mathbf{D}_G \mathbf{G}_t$, can still be larger than that of the factors for efficient prices, \mathbf{F}_t , and therefore, the factors corresponding to the leading eigenvalues in the second-step PCA may come from the microstructure noise.

Table 6 gives the simulation results when n = 390, $\bar{\tau}_G = -0.6$, and $\bar{\tau}_V = -0.6$. When we increase the sample size from 78 to 390, the REEs are smaller for "diff.total", "diff.EP", "diff.MN", and "EP". However, all three methods have larger REEs for "MN". When d = 50, the REEs are even larger than 1, but they eventually decay when d becomes larger. Table 7 provides the simulation results when n = 390, $\bar{\tau}_G = -0.4$, and $\bar{\tau}_V = -0.4$. The pattern is similar to that in Table 5, i.e., DPCA* outperforms DPCA and PCA*-VECM.

5. Application

We now apply the proposed method to 1-min and 5-min intraday returns of S&P 500 Index constituents (505 stocks in total). The data were collected from Thomson Reuters Eikon database and cover a period from 29 March 2021 to 30 June 2021. For each day, the observed prices constitute an (n + 1)-by-d matrix, \mathcal{X} , with $n \leq 78$ (for 5-min returns) or $n \leq 390$ (for 1-min returns) and $d \leq 505$. The value of n (i.e., the number of observations) and d (i.e., the number of stocks) may vary from day to day due to contemporaneous missing values at a time or suspension of trading in one day. For asynchronous missing data, we fill them using Next Observation Carried Backwards

Table 4: A when $n = 2$	verage RE] $^{78}, \bar{\tau}_G = -0$	Ξ and star $0.6, \ \overline{\tau}_V = -$	idard devi –0.6.	ation (in _F	oarenthese	s) for estin	nation of	common c	omponent	s by DPC	A^* , DPCA	r and PC≜	.*-VECM,
			DP(∂A*			DP0	CA			PCA*-1	VECM	
	diff.total	diff.EP	diff.MN	EP	MN	diff.EP	diff.MN	EP	MN	diff.EP	diff.MN	EP	MN
					R	EE unde	r Froben	ius norm					
d=50	0.209	0.261	0.283	0.284	0.690	0.266	0.339	0.359	0.731	0.306	0.319	0.278	1.006
	(0.021)	(0.111)	(0.059)	(0.092)	(0.383)	(0.116)	(0.106)	(0.134)	(0.371)	(0.080)	(0.121)	(0.113)	(0.381)
d=100	0.173	0.207	0.236	0.243	0.527	0.211	0.266	0.282	0.557	0.270	0.293	0.231	0.899
	(0.014)	(0.065)	(0.046)	(0.067)	(0.255)	(0.068)	(0.086)	(0.109)	(0.272)	(0.074)	(0.111)	(0.075)	(0.292)
d=300	0.143	0.166	0.199	0.209	0.389	0.168	0.202	0.208	0.403	0.245	0.273	0.200	0.859
	(0.010)	(0.038)	(0.041)	(0.061)	(0.173)	(0.042)	(0.075)	(0.091)	(0.216)	(0.084)	(0.117)	(0.066)	(0.354)
d=500	0.136	0.157	0.189	0.198	0.351	0.159	0.186	0.190	0.362	0.239	0.267	0.193	0.832
	(0.00)	(0.030)	(0.037)	(0.055)	(0.152)	(0.034)	(0.073)	(0.084)	(0.182)	(0.087)	(0.121)	(0.065)	(0.337)
						REE une	der MAX	r norm					
d=50	0.210	0.253	0.271	0.284	0.583	0.259	0.338	0.363	0.620	0.293	0.317	0.266	0.854
	(0.049)	(0.100)	(0.082)	(0.113)	(0.276)	(0.103)	(0.142)	(0.154)	(0.276)	(0.103)	(0.141)	(0.104)	(0.356)
d=100	0.176	0.208	0.225	0.253	0.443	0.211	0.261	0.295	0.472	0.257	0.296	0.229	0.741
	(0.034)	(0.064)	(0.061)	(0.083)	(0.180)	(0.068)	(0.113)	(0.119)	(0.202)	(0.085)	(0.122)	(0.079)	(0.290)
d=300	0.158	0.176	0.196	0.240	0.332	0.177	0.200	0.235	0.338	0.236	0.290	0.203	0.684
	(0.029)	(0.043)	(0.054)	(0.078)	(0.119)	(0.046)	(0.081)	(0.099)	(0.162)	(0.088)	(0.122)	(0.070)	(0.325)
d=500	0.157	0.170	0.190	0.235	0.305	0.172	0.188	0.219	0.305	0.232	0.288	0.200	0.659
	(0.027)	(0.039)	(0.050)	(0.074)	(0.107)	(0.042)	(0.076)	(0.085)	(0.137)	(0.086)	(0.122)	(0.065)	(0.309)
					I	REE und	er spectr	al norm					
d=50	0.161	0.257	0.235	0.230	0.660	0.261	0.315	0.302	0.677	0.259	0.258	0.274	0.983
	(0.022)	(0.126)	(0.067)	(0.079)	(0.389)	(0.130)	(0.131)	(0.122)	(0.375)	(0.087)	(0.099)	(0.129)	(0.379)
d=100	0.126	0.201	0.194	0.203	0.495	0.204	0.235	0.240	0.513	0.230	0.239	0.226	0.874
	(0.013)	(0.072)	(0.057)	(0.063)	(0.262)	(0.076)	(0.103)	(0.099)	(0.272)	(0.082)	(0.094)	(0.084)	(0.294)
d=300	0.115	0.161	0.166	0.188	0.356	0.163	0.167	0.184	0.368	0.215	0.228	0.197	0.834
	(0.010)	(0.041)	(0.053)	(0.060)	(0.176)	(0.045)	(0.082)	(0.078)	(0.212)	(0.088)	(0.095)	(0.074)	(0.353)
d=500	0.113	0.153	0.157	0.182	0.317	0.155	0.154	0.169	0.329	0.210	0.224	0.191	0.806
	(0.010)	(0.033)	(0.046)	(0.054)	(0.155)	(0.037)	(0.078)	(0.072)	(0.180)	(0.089)	(0.096)	(0.072)	(0.340)

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Table 5: A when $n = 7$	verage RE ¹ 78, $\overline{\tau}_G = -1$	E and star $0.4, \ \overline{\tau}_V =$	ndard devi. —0.4.	ation (in I	parenthese	s) for esti	mation of e	common c	omponent	s by DPC	A*, DPC/	A and PC/	*-VECM
			DP(A*			DPC	CA			PCA*-	VECM	
	diff.total	diff.EP	diff.MN	EP	MN	diff.EP	diff.MN	EP	MN	diff.EP	diff.MN	EP	MN
					R	EE unde	r Froben	ius norm					
d=50	0.185	0.306	0.391	0.151	0.360	0.514	1.562	0.778	1.088	0.494	0.216	0.403	0.847
	(0.021)	(0.107)	(0.058)	(0.039)	(0.201)	(0.250)	(0.595)	(0.280)	(0.427)	(0.153)	(0.096)	(0.160)	(0.323)
d=100	0.151	0.476	1.478	0.742	1.049	0.257	0.322	0.131	0.295	0.434	0.199	0.364	0.803
	(0.014)	(0.235)	(0.652)	(0.313)	(0.484)	(0.069)	(0.050)	(0.031)	(0.140)	(0.144)	(0.085)	(0.155)	(0.282)
d=300	0.123	0.457	1.444	0.722	1.030	0.221	0.269	0.116	0.251	0.396	0.189	0.337	0.802
	(0.010)	(0.241)	(0.706)	(0.339)	(0.508)	(0.044)	(0.042)	(0.027)	(0.126)	(0.165)	(0.093)	(0.143)	(0.292)
d=500	0.117	0.214	0.256	0.112	0.246	0.450	1.435	0.717	1.042	0.387	0.186	0.333	0.789
	(0.010)	(0.041)	(0.043)	(0.026)	(0.134)	(0.241)	(0.723)	(0.346)	(0.524)	(0.179)	(0.098)	(0.158)	(0.322)
						REE un	der MAX	norm					
d=50	0.196	0.316	0.398	0.158	0.310	0.527	1.644	0.795	0.964	0.487	0.216	0.402	0.683
	(0.070)	(0.115)	(0.106)	(0.084)	(0.172)	(0.239)	(0.708)	(0.280)	(0.379)	(0.193)	(0.111)	(0.169)	(0.299)
d=100	0.157	0.489	1.548	0.759	0.929	0.273	0.322	0.139	0.247	0.422	0.196	0.367	0.636
	(0.041)	(0.231)	(0.753)	(0.303)	(0.431)	(0.080)	(0.080)	(0.043)	(0.102)	(0.167)	(0.087)	(0.159)	(0.267)
d=300	0.138	0.469	1.485	0.739	0.896	0.244	0.276	0.138	0.212	0.392	0.194	0.342	0.624
	(0.030)	(0.231)	(0.783)	(0.320)	(0.423)	(0.060)	(0.064)	(0.038)	(0.084)	(0.178)	(0.094)	(0.144)	(0.271)
d=500	0.134	0.241	0.269	0.138	0.208	0.465	1.487	0.734	0.909	0.384	0.192	0.338	0.616
	(0.028)	(0.059)	(0.062)	(0.038)	(0.087)	(0.234)	(0.785)	(0.325)	(0.448)	(0.179)	(0.096)	(0.150)	(0.298)
					Ĥ	SEE und	er operat	or norm					
d=50	0.119	0.297	0.317	0.122	0.346	0.485	1.775	0.739	0.950	0.444	0.184	0.400	0.839
	(0.020)	(0.119)	(0.068)	(0.036)	(0.203)	(0.276)	(0.711)	(0.271)	(0.384)	(0.191)	(0.090)	(0.179)	(0.323)
d=100	0.089	0.448	1.683	0.705	0.914	0.250	0.254	0.106	0.280	0.400	0.167	0.363	0.796
	(0.012)	(0.256)	(0.784)	(0.300)	(0.438)	(0.076)	(0.061)	(0.027)	(0.144)	(0.183)	(0.080)	(0.175)	(0.284)
d=300	0.078	0.432	1.644	0.687	0.898	0.215	0.219	0.099	0.235	0.386	0.159	0.341	0.795
	(0.008)	(0.264)	(0.842)	(0.324)	(0.455)	(0.049)	(0.052)	(0.023)	(0.130)	(0.205)	(0.088)	(0.162)	(0.291)
d=500	0.077	0.210	0.212	0.097	0.229	0.425	1.635	0.681	0.910	0.382	0.157	0.339	0.782
	(0.008)	(0.047)	(0.051)	(0.021)	(0.138)	(0.264)	(0.865)	(0.330)	(0.473)	(0.219)	(0.092)	(0.178)	(0.323)

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• 5: Average REE and standard deviation (in pa	$n = 78$, $\bar{\tau}_{C} = -0.4$, $\bar{\tau}_{V} = -0.4$.
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able 5: Average REE and standard deviation (in pa	then $n = 78$, $\overline{\tau}_{C} = -0.4$, $\overline{\tau}_{V} = -0.4$.

$= 390, \bar{\tau}_G = -$	$-0.6, \bar{\tau}_V =$	0.0										
		DP(3A*			DP(CA			PCA*-	VECM	
ff.total	diff.EP	diff.MN	EP	MN	diff.EP	diff.MN	EP	MN	diff.EP	diff.MN	EP	MN
				R	EE unde	r Frobeni	ius norm					
0.177	0.214	0.222	0.257	1.524	0.214	0.258	0.317	1.500	0.197	0.213	0.200	1.548
(0.014)	(0.102)	(0.050)	(0.080)	(0.933)	(0.102)	(0.066)	(0.102)	(0.893)	(0.022)	(0.032)	(0.097)	(0.609)
0.131	0.156	0.166	0.193	1.117	0.156	0.184	0.223	1.113	0.153	0.170	0.150	1.238
(0.008)	(0.069)	(0.034)	(0.058)	(0.663)	(0.069)	(0.043)	(0.069)	(0.654)	(0.019)	(0.031)	(0.067)	(0.425)
0.089	0.100	0.112	0.133	0.630	0.100	0.116	0.140	0.628	0.115	0.138	0.102	0.943
(0.004)	(0.034)	(0.020)	(0.035)	(0.362)	(0.034)	(0.022)	(0.037)	(0.359)	(0.021)	(0.035)	(0.034)	(0.288)
0.078	0.087	0.098	0.117	0.494	0.087	0.097	0.114	0.491	0.106	0.131	0.092	0.863
(0.003)	(0.026)	(0.018)	(0.030)	(0.285)	(0.026)	(0.016)	(0.026)	(0.284)	(0.021)	(0.035)	(0.026)	(0.253)
					REE un	der MAX	norm					
0.173	0.204	0.210	0.243	1.147	0.205	0.255	0.309	1.139	0.188	0.203	0.192	1.223
(0.031)	(0.094)	(0.062)	(0.086)	(0.520)	(0.094)	(0.090)	(0.124)	(0.514)	(0.037)	(0.046)	(0.088)	(0.416)
0.125	0.147	0.154	0.182	0.833	0.148	0.178	0.215	0.831	0.141	0.162	0.142	0.947
(0.020)	(0.060)	(0.041)	(0.063)	(0.375)	(0.060)	(0.058)	(0.081)	(0.368)	(0.030)	(0.041)	(0.058)	(0.302)
0.084	0.097	0.104	0.134	0.470	0.098	0.109	0.142	0.468	0.106	0.139	0.099	0.687
(0.012)	(0.029)	(0.027)	(0.045)	(0.202)	(0.029)	(0.032)	(0.047)	(0.201)	(0.030)	(0.046)	(0.030)	(0.210)
0.077	0.087	0.092	0.125	0.369	0.087	0.090	0.121	0.366	0.100	0.136	0.091	0.629
(0.011)	(0.025)	(0.024)	(0.041)	(0.155)	(0.025)	(0.022)	(0.033)	(0.152)	(0.028)	(0.043)	(0.025)	(0.201)
				н	EE und	er operat	or norm					
0.151	0.214	0.191	0.214	1.516	0.214	0.244	0.262	1.477	0.170	0.183	0.201	1.543
(0.016)	(0.116)	(0.050)	(0.057)	(0.934)	(0.116)	(0.080)	(0.091)	(0.900)	(0.026)	(0.022)	(0.113)	(0.609)
0.107	0.155	0.141	0.155	1.109	0.155	0.169	0.182	1.098	0.129	0.138	0.149	1.233
(0.010)	(0.079)	(0.036)	(0.044)	(0.665)	(0.079)	(0.053)	(0.062)	(0.659)	(0.024)	(0.023)	(0.078)	(0.425)
0.062	0.098	0.093	0.107	0.620	0.098	0.099	0.114	0.617	0.095	0.111	0.100	0.937
(0.004)	(0.040)	(0.028)	(0.034)	(0.365)	(0.040)	(0.031)	(0.036)	(0.362)	(0.028)	(0.034)	(0.039)	(0.289)
0.053	0.085	0.080	0.098	0.484	0.085	0.077	0.094	0.481	0.089	0.108	0.090	0.857
(0.003)	(0.030)	(0.026)	(0.031)	(0.288)	(0.030)	(0.022)	(0.026)	(0.288)	(0.028)	(0.034)	(0.029)	(0.254)

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Table 7: A when $n = 1$	verage RE. 390 $\bar{\tau}_{cr} = -$	E and star -0.4 , $\overline{\tau}_V =$	ndard devi —0 4	ation (in p	oarenthese	es) for esti	mation of	common c	omponent	s by DPC	A*, DPC/	A and PC/	A*-VECM,
	D. 622	4. (DPC	3A*			DP	CA			PCA*-	VECM	
	diff.total	diff.EP	diff.MN	EP	MN	diff.EP	diff.MN	EP	MN	diff.EP	diff.MN	EP	MN
					В	EE unde	r Froben	ius norm					
d=50	0.152	0.234	0.342	0.101	0.480	0.258	0.825	0.339	0.637	0.358	0.112	0.257	0.864
	(0.014)	(0.096)	(0.035)	(0.024)	(0.287)	(0.117)	(0.414)	(0.186)	(0.328)	(0.039)	(0.026)	(0.099)	(0.292)
d=100	0.111	0.177	0.250	0.074	0.353	0.196	0.593	0.242	0.477	0.278	0.090	0.211	0.805
	(0.001)	(0.064)	(0.021)	(0.012)	(0.204)	(0.088)	(0.357)	(0.160)	(0.273)	(0.040)	(0.022)	(0.081)	(0.273)
d=300	0.074	0.126	0.167	0.054	0.211	0.139	0.392	0.161	0.310	0.209	0.076	0.172	0.764
	(0.003)	(0.032)	(0.013)	(0.007)	(0.114)	(0.053)	(0.308)	(0.135)	(0.209)	(0.042)	(0.022)	(0.059)	(0.257)
d=500	0.064	0.115	0.145	0.048	0.176	0.128	0.332	0.136	0.264	0.192	0.073	0.163	0.752
	(0.003)	(0.024)	(0.012)	(0.006)	(0.095)	(0.053)	(0.302)	(0.132)	(0.211)	(0.045)	(0.022)	(0.055)	(0.266)
						REE un	der MAX	r norm					
d=50	0.162	0.243	0.356	0.106	0.382	0.278	0.943	0.379	0.533	0.363	0.117	0.259	0.648
	(0.070)	(0.105)	(0.072)	(0.069)	(0.202)	(0.128)	(0.525)	(0.219)	(0.258)	(0.070)	(0.069)	(0.103)	(0.242)
d=100	0.110	0.184	0.253	0.073	0.269	0.209	0.670	0.270	0.386	0.269	0.088	0.212	0.574
	(0.034)	(0.063)	(0.047)	(0.033)	(0.133)	(0.086)	(0.436)	(0.185)	(0.223)	(0.057)	(0.037)	(0.079)	(0.216)
d=300	0.073	0.137	0.166	0.058	0.160	0.153	0.436	0.184	0.249	0.197	0.075	0.175	0.519
	(0.014)	(0.036)	(0.030)	(0.014)	(0.062)	(0.053)	(0.356)	(0.147)	(0.168)	(0.050)	(0.023)	(0.059)	(0.197)
d=500	0.066	0.131	0.147	0.057	0.132	0.144	0.365	0.158	0.212	0.183	0.073	0.168	0.514
	(0.013)	(0.035)	(0.028)	(0.014)	(0.051)	(0.053)	(0.333)	(0.140)	(0.167)	(0.049)	(0.022)	(0.056)	(0.208)
						REE und	er operat	or norm					
d=50	0.104	0.231	0.291	0.091	0.476	0.250	0.925	0.321	0.572	0.299	0.102	0.254	0.860
	(0.015)	(0.110)	(0.037)	(0.021)	(0.286)	(0.128)	(0.502)	(0.188)	(0.309)	(0.039)	(0.022)	(0.113)	(0.291)
d=100	0.071	0.174	0.206	0.063	0.350	0.190	0.661	0.229	0.431	0.228	0.080	0.209	0.804
	(0.007)	(0.073)	(0.022)	(0.012)	(0.205)	(0.093)	(0.430)	(0.160)	(0.253)	(0.046)	(0.021)	(060.0)	(0.273)
d=300	0.041	0.122	0.127	0.041	0.207	0.134	0.430	0.153	0.280	0.178	0.064	0.171	0.763
	(0.003)	(0.036)	(0.017)	(0.007)	(0.115)	(0.054)	(0.369)	(0.134)	(0.191)	(0.058)	(0.022)	(0.067)	(0.257)
d=500	0.033	0.112	0.106	0.039	0.171	0.123	0.359	0.130	0.239	0.170	0.060	0.162	0.751
	(0.002)	(0.027)	(0.018)	(0.007)	(0.097)	(0.055)	(0.361)	(0.129)	(0.191)	(0.063)	(0.022)	(0.062)	(0.267)

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(NOCB) if the missing data are at the beginning of the series, Last Observation Carried Forward (LOCF) if they are at the end of the series, and linear interpolation otherwise.

Firstly, we determine the number of factors using IC_1 and the PANIC tests for each of the 66 trading days within the sampling period. Figure 1 illustrates that the estimated number of factors varies from day to day. Table 8 shows some summary statistics for the estimated numbers of factors. The PANIC tests with different significance levels give similar estimates of the numbers of factors for microstructure noise and efficient prices in each day, with a difference less than 1.5 on average. Since the PANIC test at 1% performs best in the simulation, it will be used as the default PANIC test hereafter, unless specifically stated otherwise.

		$1 \cdot$	-min data				
	mean	median	1st quartile	3rd quartile	\min	\max	s.d.
K_H (IC ₁)	13.045	13	12	14	8	17	1.818
K_F (PANIC 1%)	8.894	9	8	10	4	12	1.590
K_F (PANIC 5%)	8.106	8	7	9	4	12	1.656
K_F (PANIC 10%)	7.561	8	6	9	3	12	1.890
K_G (PANIC 1%)	4.152	4	3	5	1	8	1.552
K_G (PANIC 5%)	4.939	5	4	6	2	9	1.626
K_G (PANIC 10%)	5.485	5	4	6.75	2	10	1.629
		5.	-min data				
	mean	median	1st quartile	3rd quartile	\min	max	s.d.
K_H (IC ₁)	7.697	8	7	8.75	5	12	1.488
K_F (PANIC 1%)	6.758	7	6	8	3	12	1.710
K_F (PANIC 5%)	6.076	6	5	7	2	12	1.892
K_F (PANIC 10%	5.500	5	4	7	2	10	1.774
K_G (PANIC 1%)	0.939	1	0	2	0	3	0.990
K_G (PANIC 5%)	1.621	2	0	2	0	5	1.274
K_G (PANIC 10%)	2.197	2	2	3	0	6	1.268

Table 8: Summary statistics for estimated numbers of factors over the sampling period

We can see from Table 8 that the numbers of factors are larger for 1-min data than those for 5-min data. The reason might be twofold. Firstly, the magnitude of microstructure noise decreases when 1-min data are aggregated to 5-min data. Thus the factors for microstructure noise are more difficult to detect. Secondly, due to the Epps effect (as evident in Table 10), the correlation between stocks decreases as the sampling frequency increases, resulting in higher numbers of factors for both microstructure noise and efficient prices of higher frequency data.

To see how the number of factors change during the sample period, we looked at whether there is a relation between the number of factors and the following variables: the market excess return (MKT), the size factor (SMB), the value factor (HML), the Momentum factor (MOM), the short-term Reversal factor (STREV), the long-term Reversal factor (LTREV), the CBOE Volatility Index (VIX), the high low volatility on the S&P500 index (HLVOL) and the implied Roll measure of bid-ask spread (ROLL, computed as $2\sqrt{-[Cov(r_t, r_{t-1}) \wedge 0]}$ using 5-min price changes r_t and averaged over

all stocks with equal weight). The first 6 risk factors were downloaded from Kenneth R. French's data library² and the VIX and S&P500 index were downloaded from Thomson Reuters Eikon database.

We find from Table 9 that the contemporaneous correlations between the number of factors and the above-mentioned risk variables are insignificant. However, some of these variables can partially explain changes in the number of factors for efficient prices in the next day. Specifically, the highest (in absolute value) correlation is between lag-1 SMB and 1-min K_F , which is 0.329 with a p-value of 0.007. Significant correlations are also found, at 10% significant level, between 1-min K_F and HLVOL at -0.240, between 1-min K_F and MOM at 0.211, and between 1-min K_F and LTREV at 0.213. What's more, we also find a negative relation between the number of factors and VIX (or HLVOL). Figure 2 shows this relation in a time series plot: when the S&P500 high-low volatility peaks on 12 May 2021, the 1-min K_H and 1-min K_F drop in the next day. This indicates that the co-movement of stocks increases during High VIX (or HLVOL) period, confirming the old adage that diversification disappears when needed most.

Table 9: Correlation between number of factors and 9 risk variables. Panel A shows the contemporaneous correlation and Panel B shows the lagged correlation (the 9 risk variables are lagged by 1 day). P-value less than 0.1, 0.05 or 0.01 is flagged with one, two or three stars (*, **, ***), respectively.

(Panel A)	VIX	HLVOL	MKT	SMB	HML	MOM	STREV	LTREV	ROLL
$K_H(1-\min)$	-0.193	-0.121	0.043	0.185	-0.038	0.178	0.141	0.153	-0.034
$K_F(1\text{-min})$	-0.134	-0.033	0.004	0.159	-0.054	0.065	-0.010	0.094	-0.113
$K_G(1\text{-min})$	-0.089	-0.108	0.047	0.055	0.010	0.142	0.175	0.083	0.075
$K_H(5\text{-min})$	-0.008	0.123	0.042	0.107	0.115	0.086	0.090	0.205	0.083
$K_F(5\text{-min})$	-0.010	0.093	0.009	0.004	-0.061	-0.011	-0.051	-0.005	0.010
$K_G(5\text{-min})$	0.004	0.025	0.048	0.155	0.279	0.149	0.223	0.317	0.106
(Panel B)	VIX	HLVOL	MKT	SMB	HML	MOM	STREV	LTREV	ROLL
$K_H(1-\min)$	-0.206^{*}	-0.191	0.157	0.292**	0.129	0.231^{*}	-0.184	0.254^{**}	-0.132
$K_F(1\text{-min})$	-0.131	-0.240^{*}	0.144	0.329^{***}	* 0.050	0.211^{*}	-0.027	0.213^{*}	0.005
$K_G(1\text{-min})$	-0.105	0.023	0.036	0.002	0.099	0.052	-0.187	0.078	-0.159
$K_H(5\text{-min})$	0.003	-0.088	0.110	0.026	0.041	-0.084	-0.072	-0.006	0.033
$K_F(5\text{-min})$	-0.013	-0.076	0.183	-0.019	0.055	-0.085	-0.150	-0.032	-0.056
$K_G(5-\min)$	0.027	-0.003	-0.146	0.071	-0.032	0.019	0.146	0.046	0.144

For illustration purposes, we consider the results for 30 June 2021, i.e., the last day in the sample period. For 1-min data, the estimated total number of factors is 13, among which 6 are identified as factors for efficient prices by the PANIC test. Figure 3 shows the 13 estimated factors in cumulative form, i.e., $(\hat{\boldsymbol{\beta}}_{\perp}, \hat{\boldsymbol{\beta}})^{\mathsf{T}} \widehat{\boldsymbol{H}}_{s\Delta}^*$, where $(\hat{\boldsymbol{\beta}}_{\perp}, \hat{\boldsymbol{\beta}})$ is the matrix of eigenvectors of the matrix $n^{-1} \widehat{\mathscr{H}}^{*c^{\mathsf{T}}} \widehat{\mathscr{H}}^{*c}$, arranged in descending order of their corresponding eigenvalues. The first 6 factors appear to be more variable than the last 7 factors.

We can not tell from Figure 3 whether the factors for microstructure noise dominate those for efficient prices, as estimated factors have been standardised. Instead, we calculate the variance ratio

²http://mba.tuck.dartmouth.edu/pages/faculty/ken.french/data_library.html.

			Correl	ation				
	mean	median	1st quartile	3rd quartile	min	max	s.d.	
1-min	0.083	0.075	-0.025	0.179	-0.471	0.925	0.160	
5-min	0.097	0.097	-0.062	0.253	-0.674	0.960	0.229	
Absolute value of correlation								
	mean	median	1st quartile	3rd quartile	min	max	s.d.	
1-min	0.139	0.112	0.053	0.195	0.000	0.925	0.114	
5-min	0.200	0.171	0.082	0.288	0.000	0.960	0.148	

Table 10: Summary statistics for pairwise correlations between stock returns on 30 June 2021

Table 11: Summary statistics for the variance ratio of the common component for microstructure noise to that of efficient price for each stock on 30 June 2021. Note that we define both variance ratio based on $\widehat{\operatorname{Var}}(\widehat{\lambda}_{Fi}^{*^{\mathsf{T}}}\widehat{f}_{s\Delta}^{*})$, as the variance $\widehat{\operatorname{Var}}(\widehat{\lambda}_{Fi}^{*^{\mathsf{T}}}\widehat{f}_{s\Delta}^{*})$ is not meaningful.

٢	Variano	ce ratio d	of common c	components:	$\widehat{\operatorname{Var}}(\widehat{\boldsymbol{\lambda}}_{G}^{*})$	$\widehat{G}_{s\Delta}^*)/\widetilde{\mathbb{V}}$	$\widehat{\operatorname{Var}}(\widehat{oldsymbol{\lambda}}_{Fi}^{*^\intercal}\widehat{oldsymbol{f}}_{s\Delta}^*)$	
	mean	median	1st quartile	3rd quartile	min	max	s.d.	
1-min	4.116	2.541	1.302	5.451	0.030	48.788	4.875	
5-min	0.148	0.048	0.013	0.138	0.000	4.228	0.325	
$\textbf{Variance ratio of differenced common components: } \widehat{\text{Var}}(\widehat{\boldsymbol{\lambda}}_{Gi}^{*^\intercal}\widehat{\boldsymbol{g}}_{s\Delta}^*)/\widehat{\text{Var}}(\widehat{\boldsymbol{\lambda}}_{Fi}^{*^\intercal}\widehat{\boldsymbol{f}}_{s\Delta}^*)}$								
	mean	median	1st quartile	3rd quartile	min	max	s.d.	
1-min	2.633	1.604	0.812	3.492	0.017	35.849	3.229	
5-min	0.069	0.020	0.006	0.062	0.000	2.209	0.156	

of the common component for microstructure noise to that of efficient price for each stock, to take the magnitude of factor loadings into consideration. We give the summary statistics in Table 11. We can see that on average, the common component for microstructure noise dominates the common component for efficient prices at 1-min frequency, while the relation reverses at 5-min frequency. For individual stocks, however, the contribution of the common component for microstructure noise may still be small even at the 1-min frequency. To show this, we look at the decomposition of prices (cumulative returns) into three components: the common component of efficient prices, $\Lambda_F \mathbf{F}_t$ (CC.EP), the common component of microstructure noise, $\Lambda_G \mathbf{D}_G \mathbf{G}_t$ (CC.MN), and idiosyncratic errors (Residuals). We illustrate with five randomly selected stocks that have the stock ticker symbols – POOL, CHRW, AJG, CNP, and WM. Figure 4 shows the decomposition for the cumulative 1-min returns of the five stocks on 30 June 2021. The corresponding numbers of factors are $\hat{K}_F = 6$ and $\hat{K}_G = 7$. Figure 5 shows the decomposition for 5-min data with $\hat{K}_F = 5$ and $\hat{K}_G = 1$. The two figures show that the common component of the microstructure noise can explain only a small amount of the variability of the prices.

In summary, our analysis finds existence of common components for the microstructure noise of S&P 500 stocks, although their magnitude is small. The small magnitude is also consistent with the expectation that there are very few arbitrage opportunities in a frictional market.

Figure 1: Estimated number of factors for each trading day from 29 March 2021 to 30 June 2021. The y-axis represents the number of factors and the x-axis represents the dates (given in the format mmdd). The y-coordinate of the top of each grey bar gives the estimated total number of factors, \hat{K}_H , from IC₁. The length of each grey bar represents the difference between \hat{K}_H and \hat{K}_F , which is obtained from the PANIC test using 1% significance level. The length of each red bar represents the difference between \hat{K}_F 's obtained from the PANIC tests using 1% and 5% significance levels. The length of each blue bar represents the difference between \hat{K}_F 's obtained from the PANIC tests using 1% and 5% significance levels. The length of each blue bar represents the difference between \hat{K}_F 's obtained from the PANIC tests using 5% and 10% significance levels. The y-coordinate of the bottom of each blue bar gives the value of \hat{K}_F obtained from a 10% PANIC test.



Figure 2: Time series plots of $K_H(1-\min)$ and $K_F(1-\min)$, VIX, and HLVOL over the sampling period







components of microstructure noise (CC.MN), and the idiosyncratic errors (Residuals) for the stocks POOL, CHRW, AJG, CNP, and WM, with $\hat{K}_F = 6$ and $\hat{K}_G = 7$. Each row gives the decomposition for each stock, with the first diagram giving the cumulative returns, followed Figure 4: Decomposition of cumulative 1-min returns on 30 June 2021 into the common components of efficient prices (CC.EP), the common by CC.EP, CC.MN, and Residuals. The variance ratio $\widetilde{\operatorname{Var}}(\widehat{\lambda}_{G_s}^* \widehat{\Omega}_{s\Delta}^*)/\widetilde{\operatorname{Var}}(\widehat{\lambda}_{F_s}^* \widehat{\Pi}_{S\Delta}^*)$ for each stock is 0.462, 0.980, 0.125, 1.282, and 6.818, respectively.



with $\hat{K}_F = 5$ and $\hat{K}_G = 1$. Each row gives the decomposition for each stock, with the first diagram giving the cumulative returns, followed Figure 5: Decomposition of cumulative 5-min returns on 30 June 2021 into the common components of efficient prices (CC.EP), the common by CC.EP, CC.MN, and Residuals. The variance ratio $\widetilde{\operatorname{Var}}(\widehat{\lambda}_{G_s}^* \widehat{\Omega}_{s\Delta}^*)/\widetilde{\operatorname{Var}}(\widehat{\lambda}_{F_s}^* \widehat{\Pi}_{S\Delta}^*)$ for each stock is 0.183, 0.015, 0.001, 0.180, and 0.550, components of microstructure noise (CC.MN), and the idiosyncratic errors (Residuals) for the stocks POOL, CHRW, AJG, CNP, and WM respectively.



6. Conclusion and Future Work

We consider a dual-factor model for high-frequency stock prices contaminated with microstructure noise. We develop the Double Principle Component Analysis (DPCA*) to estimate the separate factor structures for efficient prices and microstructure noise. When comparing with the PCA-VECM approach, the DPCA* approach is free from the need to impose strong parametric assumptions on the microstructure noise and applies instead to a broad class of stationary processes. The estimators are proven to be consistent and perform well in simulations. The empirical analysis of intraday returns of S&P 500 constituents provides some evidence of co-movement in the microstructure noise, apart from co-movement of prices caused by common systematic risk factors.

Identifying and separating out common components of microstructure noise from the prices are very useful for the study of properties of the microstructure noise and efficient price processes. For example, once the common components for microstructure noise are separated out, the common components for efficient prices are no longer contaminated by microstructure noise and hence, can be used to obtain a more accurate estimate of the common part of realized volatility. For the idiosyncratic part of realized volatility, one can use the estimated idiosyncratic errors and apply the pre-averaging method of Jacod et al. (2009). Adding these two parts together, we get an estimator of the realized volatility matrix. We may introduce sparsity or block structure into idiosyncratic components like Dai et al. (2019) and Aït-Sahalia and Xiu (2017), respectively. However, since our main interests are the identification of common factors, we avoid introducing these structures and leave the estimation of the realized volatility matrix to the future work.

The estimated common factors and loadings for microstructure noise provide useful tools for portfolio management. With such estimates, portfolio managers can construct a new factor mimicking portfolio which is only exposed to the factors of microstructure noise. Such a portfolio can be used to hedge risks from microstructure noise. Since the portfolio return is usually stationary, one can apply the mean-reverting strategy to earn profits from the portfolio, once its volatility is large enough to cover the cost. Even if its volatility is small, one can still time the market according to it, e.g., when adjusting the position of a portfolio, to lower the cost.

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Appendix A: Technical Proofs

In the subsequent proofs, we often make use of the following Weyl's inequality, for two $n \times n$ symmetric matrices \mathbf{M}_1 and \mathbf{M}_2 , with eigenvalues $\mu_j(\mathbf{M}_1)$ and $\mu_j(\mathbf{M}_2)$:

$$|\mu_j(\mathbf{M}_1) - \mu_j(\mathbf{M}_2)| \le \|\mathbf{M}_1 - \mathbf{M}_2\|_{sp},$$
 (A.1)

for j = 1, ..., n. If \mathbf{M}_1 and \mathbf{M}_2 are invertible and $\|\mathbf{M}_1 - \mathbf{M}_2\|_{sp} \|\mathbf{M}_2^{-1}\|_{sp} < 1$, we have

$$\begin{split} \|\mathbf{M}_{1}^{-1} - \mathbf{M}_{2}^{-1}\|_{\mathsf{sp}} &\leq \|\mathbf{M}_{1}^{-1}\|_{\mathsf{sp}} \|\mathbf{M}_{1} - \mathbf{M}_{2}\|_{\mathsf{sp}} \|\mathbf{M}_{2}^{-1}\|_{\mathsf{sp}} \\ &\leq \|\mathbf{M}_{1}^{-1} - \mathbf{M}_{2}^{-1}\|_{\mathsf{sp}} \|\mathbf{M}_{1} - \mathbf{M}_{2}\|_{\mathsf{sp}} \|\mathbf{M}_{2}^{-1}\|_{\mathsf{sp}} \\ &\leq \frac{\|\mathbf{M}_{2}^{-1}\|_{\mathsf{sp}} \|\mathbf{M}_{1} - \mathbf{M}_{2}\|_{\mathsf{sp}} \|\mathbf{M}_{2}^{-1}\|_{\mathsf{sp}}}{1 - \|\mathbf{M}_{1} - \mathbf{M}_{2}\|_{\mathsf{sp}} \|\mathbf{M}_{2}^{-1}\|_{\mathsf{sp}}}. \end{split}$$
(A.2)

Note that the max norm is not submultiplicative, but we can use $\|\mathbf{M}_1\mathbf{M}_2\|_{\mathsf{MAX}} \le \|\mathbf{M}_1\|_{\infty} \|\mathbf{M}_2\|_{\mathsf{MAX}}$ or $\|\mathbf{M}_1\mathbf{M}_2\|_{\mathsf{MAX}} \le \|\mathbf{M}_1\|_{\mathsf{MAX}} \|\mathbf{M}_2\|_1$.

Lemma A.1. Under Assumptions 1–4, 5^{*}, 6 and 7, we have, for each $1 \le j \le K_H$, (i)

$$\xi_j - \frac{\mathbf{b}_j}{\|\mathbf{b}_j\|_2} = O\left(\frac{m_{U,d} + n^{2\bar{\tau}_V^*} m_{v,d}}{dn^{2\bar{\tau}_G^{*-}}}\right),\tag{A.3}$$

where ξ_j is the eigenvector of $\Lambda_H \mathbf{D}_H \boldsymbol{\Sigma}_h \mathbf{D}_H \Lambda_H^{\mathsf{T}}$ corresponding to the *j*th largest eigenvalue. (*ii*) If, in addition, Assumption 8(*i*) holds, then $\mu_{K_H}(\boldsymbol{x}^{\mathsf{T}}\boldsymbol{x}) \geq C dn^{2\underline{\tau}_G^{*-}}$. *Proof.* Recall that

$$\boldsymbol{\Sigma}_x = \boldsymbol{\Lambda}_H \mathbf{D}_H \boldsymbol{\Sigma}_h \mathbf{D}_H \boldsymbol{\Lambda}_H^{\mathsf{T}} + \boldsymbol{\Sigma}_w.$$

Let $\mathbf{B} = \mathbf{\Lambda}_H \mathbf{D}_H \mathbf{\Sigma}_h^{1/2} \mathbf{Q} = (\mathbf{b}_1, \cdots, \mathbf{b}_{K_H})$ with $\|\mathbf{b}_j\|_2$'s sorted in a descending order, where \mathbf{Q} is an orthogonal matrix such that $\mathbf{Q}^{\mathsf{T}} \mathbf{\Sigma}_h^{1/2} \mathbf{D}_H \mathbf{\Lambda}_H^{\mathsf{T}} \mathbf{\Lambda}_H \mathbf{D}_H \mathbf{\Sigma}_h^{1/2} \mathbf{Q}$ is a diagonal matrix. Then $\|\mathbf{b}_j\|_2^2$, $1 \leq j \leq K_H$, are the non-zero eigenvalues of $\mathbf{B}\mathbf{B}^{\mathsf{T}} = \mathbf{\Lambda}_H \mathbf{D}_H \mathbf{\Sigma}_h \mathbf{D}_H \mathbf{\Lambda}_H^{\mathsf{T}}$ and also the eigenvalues of $\mathbf{B}^{\mathsf{T}} \mathbf{B} = \mathbf{\Sigma}_h^{1/2} \mathbf{D}_H \mathbf{\Lambda}_H^{\mathsf{T}} \mathbf{\Lambda}_H \mathbf{D}_H \mathbf{\Sigma}_h^{1/2}$. Therefore,

$$\|\boldsymbol{b}_{j}\|_{2}^{2} \leq \|\boldsymbol{b}_{1}\|_{2}^{2} = \|\boldsymbol{\Sigma}_{h}^{1/2}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{^{\mathsf{T}}}\boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\boldsymbol{\Sigma}_{h}^{1/2}\|_{\mathsf{sp}} \leq \|\mathbf{D}_{H}\|_{\mathsf{sp}}^{2}\|\boldsymbol{\Sigma}_{h}\|_{\mathsf{sp}} \cdot \|\boldsymbol{\Lambda}_{H}^{^{\mathsf{T}}}\boldsymbol{\Lambda}_{H}\|_{\mathsf{sp}} = O(dn^{2\bar{\tau}_{G}^{*+}}),$$

where the last equality holds by Assumptions 2 and 4. On the other hand,

$$\|\boldsymbol{b}_{K_{H}}\|_{2}^{2} = \mu_{K_{H}}(\boldsymbol{\Sigma}_{h}^{1/2}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{^{\mathsf{T}}}\boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\boldsymbol{\Sigma}_{h}^{1/2})$$

$$\geq \mu_{K_{H}}(\boldsymbol{\Sigma}_{h})\mu_{K_{H}}^{2}(\mathbf{D}_{H})\mu_{K_{H}}(\boldsymbol{\Lambda}_{H}^{^{\mathsf{T}}}\boldsymbol{\Lambda}_{H}) \geq Cdn^{2\underline{\tau}_{G}^{*-}}.$$
(A.4)

By the Sin theta theorem in Davis and Kahan (1970) (see Yu et al. (2015) for a statistician-friendly version), we have

$$\left|\xi_{j} - \frac{\boldsymbol{b}_{j}}{\|\boldsymbol{b}_{j}\|_{2}}\right| \leq \frac{\sqrt{2}\|\boldsymbol{\Sigma}_{w}\|_{\text{sp}}}{\min\left\{|\mu_{j-1}(\boldsymbol{\Sigma}_{x}) - \|\boldsymbol{b}_{j}\|_{2}^{2}|, |\mu_{j+1}(\boldsymbol{\Sigma}_{x}) - \|\boldsymbol{b}_{j}\|_{2}^{2}|\right\}},\tag{A.5}$$

with the convention $\mu_0(\cdot) = \infty$. By Weyl's inequality and triangle inequality, we have

$$|\mu_j(\boldsymbol{\Sigma}_x) - \|\boldsymbol{b}_j\|_2^2| \le \|\boldsymbol{\Sigma}_w\|_{\mathsf{sp}} = O(m_{w,nd}), \tag{A.6}$$

for $1 \leq j \leq K_H$ and

$$0 < \mu_j(\boldsymbol{\Sigma}_x) \le \|\boldsymbol{\Sigma}_w\|_{\mathsf{sp}} = O(m_{w,nd}),$$

for $K_H + 1 \leq j \leq d$, where $m_{w,nd} = m_{U,d} + n^{2\bar{\tau}_V^*} m_{v,d}$. Using the triangle inequality, (A.6) and Assumption 7, we have

$$\begin{aligned} \left| \mu_{j-1}(\boldsymbol{\Sigma}_{x}) - \|\boldsymbol{b}_{j}\|_{2}^{2} \right| &\geq \left| \|\boldsymbol{b}_{j-1}\|_{2}^{2} - \|\boldsymbol{b}_{j}\|_{2}^{2} \right| - \left| \mu_{j-1}(\boldsymbol{\Sigma}_{x}) - \|\boldsymbol{b}_{j-1}\|_{2}^{2} \right| \\ &\geq C dn^{2\underline{\tau}_{G}^{*-}}, \end{aligned}$$

for $1 \leq j \leq K_H$, as $m_{w,nd} = o(dn^{2\underline{\tau}_G^*})$ under Assumption 5*. Similarly, $|\mu_{j+1}(\Sigma_x) - ||\mathbf{b}_j||_2^2| \geq Cdn^{2\underline{\tau}_G^*}$, when $1 \leq j \leq K_H - 1$. Therefore, by (A.5), we can prove (i). For part (ii), by Weyl's inequality,

 $|\mu_{K_H}(x^{^{\intercal}}x)-\mu_{K_H}(\Sigma_x)|\leq \|x^{^{\intercal}}x-\Sigma_x\|_{\mathsf{sp}}.$

Then using $\mu_{K_H}(\boldsymbol{\Sigma}_x) \geq \|\boldsymbol{b}_{K_H}\|_2^2$, we have

$$\begin{split} \mu_{K_H}(x^{\mathsf{T}}x) &\geq & \mu_{K_H}(\boldsymbol{\Sigma}_x) - \|x^{\mathsf{T}}x - \boldsymbol{\Sigma}_x\|_{\mathsf{sp}} \\ &\geq & \|\boldsymbol{b}_{K_H}\|_2^2 - \|x^{\mathsf{T}}x - \boldsymbol{\Lambda}_H \mathbf{D}_H \boldsymbol{\lambda}^{\mathsf{T}} \boldsymbol{\hbar} \mathbf{D}_H \boldsymbol{\Lambda}_H^{\mathsf{T}}\|_{\mathsf{sp}} - \|\boldsymbol{\Sigma}_w\|_{\mathsf{sp}}, \\ & & - \|\boldsymbol{\Lambda}_H \mathbf{D}_H \boldsymbol{\lambda}^{\mathsf{T}} \boldsymbol{\hbar} \mathbf{D}_H \boldsymbol{\Lambda}_H^{\mathsf{T}} - \boldsymbol{\Lambda}_H \mathbf{D}_H \boldsymbol{\Sigma}_h \mathbf{D}_H \boldsymbol{\Lambda}_H^{\mathsf{T}}\|_{\mathsf{sp}} \end{split}$$

where $\Sigma_w = \Sigma_U + n \mathbf{D}_V \Sigma_v \mathbf{D}_V$. Therefore by (A.4), (A.6) and Assumption 5*, we only need to show that

$$\|\boldsymbol{x}^{\mathsf{T}}\boldsymbol{x} - \boldsymbol{\Lambda}_{H}\boldsymbol{\mathrm{D}}_{H}\boldsymbol{\hbar}^{\mathsf{T}}\boldsymbol{\hbar}\boldsymbol{\mathrm{D}}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\|_{\mathsf{sp}} = o_{P}(dn^{2\underline{\tau}_{G}^{*-}}),\tag{A.7}$$

and

$$\|\boldsymbol{\Lambda}_{H}\boldsymbol{\mathrm{D}}_{H}\boldsymbol{\hbar}^{\mathsf{T}}\boldsymbol{\hbar}\boldsymbol{\mathrm{D}}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}-\boldsymbol{\Lambda}_{H}\boldsymbol{\mathrm{D}}_{H}\boldsymbol{\Sigma}_{h}\boldsymbol{\mathrm{D}}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\|_{\mathsf{sp}}=o_{P}(dn^{2\underline{\tau}_{G}^{*-}}).$$
(A.8)

As for (A.7), using Lemma B.4, we have

$$\begin{aligned} \|x^{\mathsf{T}}x - \mathbf{\Lambda}_{H}\mathbf{D}_{H}\boldsymbol{\Lambda}^{\mathsf{T}}_{H}\|_{\mathsf{sp}} &\leq 2\|\mathbf{\Lambda}_{H}\|_{\mathsf{sp}}\|\mathbf{D}_{H}\boldsymbol{\Lambda}^{\mathsf{T}}w\|_{\mathsf{sp}} + \|w^{\mathsf{T}}w - \mathbf{\Sigma}_{w}\|_{\mathsf{sp}} + \|\mathbf{\Sigma}_{w}\|_{\mathsf{sp}} \\ &= O_{P}(d(\log d/n)^{1/2} \cdot n^{\bar{\tau}_{V}^{*+} + \bar{\tau}_{G}^{*+}}) \\ &+ O_{P}(d(\log d/n)^{1/2} \cdot n^{2\bar{\tau}_{V}^{*+}}) + O(m_{w,nd}), \end{aligned}$$
(A.9)

since $\|\mathbf{\Lambda}_H\|_{\mathsf{sp}} = O(d^{1/2}), \|\mathbf{D}_H \hbar^{\mathsf{T}} w\|_{\mathsf{sp}} \leq d^{1/2} \|\mathbf{D}_H \hbar^{\mathsf{T}} w\|_1$ and $\|w^{\mathsf{T}} w - \mathbf{\Sigma}_w\|_{\mathsf{sp}} \leq d\|w^{\mathsf{T}} w - \mathbf{\Sigma}_w\|_{\mathsf{MAX}}$. Then, under Assumptions 5* and 8(i),

$$\|\boldsymbol{x}^{\mathsf{T}}\boldsymbol{x} - \boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\boldsymbol{\hbar}^{\mathsf{T}}\boldsymbol{\hbar}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\|_{\mathsf{sp}} = o_{P}(dn^{2\underline{\tau}_{G}^{*-}}),$$

if $n^{1+4\underline{\tau}_{G}^{*-}-2\overline{\tau}_{V}^{*+}-2(\overline{\tau}_{G}^{*+}\vee\overline{\tau}_{V}^{*+})}/(\log d) \to \infty$ and $m_{w,nd} = o(dn^{2\underline{\tau}_{G}^{*-}})$. As for (A.8), by Assumptions 4 and 8(i), we have

$$\begin{split} &\|\boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\boldsymbol{\hbar}^{\mathsf{T}}\boldsymbol{\hbar}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}-\boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\boldsymbol{\Sigma}_{h}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\|_{\mathsf{sp}} \\ &\leq C_{\Lambda}d\|\mathbf{D}_{H}\boldsymbol{\hbar}^{\mathsf{T}}\boldsymbol{\hbar}\mathbf{D}_{H}-\mathbf{D}_{H}\boldsymbol{\Sigma}_{h}\mathbf{D}_{H}\|_{\mathsf{sp}} \\ &\leq K_{H}C_{\Lambda}d\cdot\max\left\{\|\boldsymbol{\ell}^{\mathsf{T}}\boldsymbol{\ell}-\boldsymbol{\Sigma}_{F}\|_{\mathsf{MAX}},\|\mathbf{D}_{G}\boldsymbol{\mathscr{G}}^{\mathsf{T}}\boldsymbol{\mathscr{G}}\mathbf{D}_{G}-\mathbf{D}_{G}\boldsymbol{\Sigma}_{g}\mathbf{D}_{G}\|_{\mathsf{MAX}},\|\boldsymbol{\ell}^{\mathsf{T}}\boldsymbol{\mathscr{G}}\mathbf{D}_{G}\|_{\mathsf{MAX}}\right\} \\ &= O_{P}(dn^{2\bar{\tau}_{G}^{*+}}(\log d/n)^{1/2}) = o_{P}(dn^{2\underline{\tau}_{G}^{*-}}). \end{split}$$

where the last line follows from Lemmas B.1(ii), B.2(ii) and B.3(iv). Hence we complete the proof.

Lemma A.2. Suppose that Assumptions 1–8 are satisfied. We have (i)

$$\left\|\widehat{\mathbf{\Lambda}}_{H} - \mathbf{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R}\right\|_{\mathsf{MAX}} = O_{P}\left(n^{-2\underline{\tau}_{G}^{*-}} \cdot a_{nd}\right) = o_{P}(1), \tag{A.10}$$

where

$$a_{nd} = (\log d)^{1/2} \frac{n^{\bar{\tau}_V^{*+} + \bar{\tau}_G^{*+} \vee \bar{\tau}_V^{*+}}}{n^{1/2}} + \frac{m_{w,nd}}{d^{1/2}},$$

in which $\bar{\tau}_G^{*+} = (1/2 + \bar{\tau}_G)_+, \ \underline{\tau}_G^{*-} = (1/2 + \underline{\tau}_G)_-, \ \bar{\tau}_V^{*+} = (1/2 + \bar{\tau}_V)_+, \ m_{w,nd} = m_{U,d} + n^{2\bar{\tau}_V^*} m_{v,d}, \ and \ \bar{\tau}_V^* = (1/2 + \bar{\tau}_V).$ (ii)

$$\|\mathbf{D}_{H}\mathbf{R}\|_{sp} = O_{P}(1), \quad \|(\mathbf{D}_{H}\mathbf{R})^{-1}\|_{sp} = O_{P}(1),$$
 (A.11)

and

$$d^{-1/2} \|\mathbf{R}\widehat{\mathbf{D}}_{x,K_H}^{1/2}\|_{\mathsf{sp}} = O_P(1).$$
(A.12)

(iii)

$$\left\|\widehat{\boldsymbol{\lambda}}^{\mathsf{T}} - (\mathbf{D}_{H}\mathbf{R})^{-1}\mathbf{D}_{H}\boldsymbol{\lambda}^{\mathsf{T}}\right\|_{\mathsf{MAX}} \leq \left\|\widehat{\boldsymbol{\lambda}}^{\mathsf{T}} - (\mathbf{D}_{H}\mathbf{R})^{-1}\mathbf{D}_{H}\boldsymbol{\lambda}^{\mathsf{T}}\right\|_{\mathsf{sp}} = O_{P}\left(n^{-\underline{\tau}_{G}^{*-}} \cdot a_{nd}\right),\tag{A.13}$$

where the rotation matrix \mathbf{R} is defined as

$$\mathbf{R} = \mathbf{\hbar}^{\mathsf{T}} \mathbf{\hbar} \mathbf{D}_{H} \mathbf{\Lambda}_{H}^{\mathsf{T}} \widehat{\mathbf{\Lambda}}_{H} \widehat{\mathbf{D}}_{x, K_{H}}^{-1}, \qquad (A.14)$$

in which $\widehat{\mathbf{D}}_{x,K_H}$ is a $K_H \times K_H$ diagonal matrix with the diagonal elements being the first K_H largest eigenvalues of $x^{\mathsf{T}}x$ arranged in descending order.

PROOF OF LEMMA A.2. By the definition of PCA estimation, we may show that

$$egin{aligned} & \left(\widehat{\mathbf{\Lambda}}_{H}-\mathbf{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R}
ight)\widehat{\mathbf{D}}_{x,K_{H}} = (x^{^{\intercal}}x-\mathbf{\Lambda}_{H}\mathbf{D}_{H}{}^{\intercal}{}^{\intercal}{}^{\hslash}\mathbf{D}_{H}\mathbf{\Lambda}_{H}^{^{\intercal}})\widehat{\mathbf{\Lambda}}_{H} \ & = \mathbf{\Lambda}_{H}\mathbf{D}_{H}{}^{\hslash}{}^{\intercal}{}^{w}\widehat{\mathbf{\Lambda}}_{H} + w^{^{\intercal}}{}^{\hslash}\mathbf{D}_{H}\mathbf{\Lambda}_{H}^{^{\intercal}}\widehat{\mathbf{\Lambda}}_{H} \ & + (w^{^{\intercal}}w-\mathbf{\Sigma}_{w})\widehat{\mathbf{\Lambda}}_{H} + \mathbf{\Sigma}_{w}\widehat{\mathbf{\Lambda}}_{H}. \end{aligned}$$

For the first term on the right hand side of the second equality, we have, by similar argument to (A.9),

$$\begin{split} \|x^{\mathsf{T}}x - \mathbf{\Lambda}_{H}\mathbf{D}_{H}\boldsymbol{\hbar}^{\mathsf{T}}\boldsymbol{\hbar}\mathbf{D}_{H}\mathbf{\Lambda}_{H}^{\mathsf{T}}\|_{\mathsf{MAX}} &\leq 2\|\mathbf{\Lambda}_{H}\|_{\mathsf{MAX}}\|\mathbf{D}_{H}\boldsymbol{\hbar}^{\mathsf{T}}\boldsymbol{w}\|_{1} + \|\boldsymbol{w}^{\mathsf{T}}\boldsymbol{w} - \boldsymbol{\Sigma}_{w}\|_{\mathsf{MAX}} + \|\boldsymbol{\Sigma}_{w}\|_{\mathsf{MAX}} \\ &= O_{P}((\log d/n)^{1/2} \cdot n^{\bar{\tau}_{V}^{*+} + \bar{\tau}_{G}^{*+}}) \\ &+ O_{P}((\log d/n)^{1/2} \cdot n^{2\bar{\tau}_{V}^{*+}}) + O(m_{w,nd}). \end{split}$$

For the remaining terms, we have

$$\begin{split} \|\boldsymbol{\Lambda}_{H} \mathbf{D}_{H} \boldsymbol{\hbar}^{\mathsf{T}} \boldsymbol{w} \widehat{\boldsymbol{\Lambda}}_{H}\|_{\mathsf{MAX}} &\leq \|\boldsymbol{\Lambda}_{H}\|_{\mathsf{MAX}} \|\mathbf{D}_{H} \boldsymbol{\hbar}^{\mathsf{T}} \boldsymbol{w}\|_{1} \|\widehat{\boldsymbol{\Lambda}}_{H}\|_{1} = O_{P}(d(\log d/n)^{1/2} \cdot n^{\bar{\tau}_{V}^{*+} + \bar{\tau}_{G}^{*+}}), \\ \|\boldsymbol{w}^{\mathsf{T}} \boldsymbol{\hbar} \mathbf{D}_{H} \boldsymbol{\Lambda}_{H}^{\mathsf{T}} \widehat{\boldsymbol{\Lambda}}_{H}\|_{\mathsf{MAX}} &\leq \|\boldsymbol{w}^{\mathsf{T}} \boldsymbol{\hbar} \mathbf{D}_{H}\|_{\mathsf{MAX}} \|\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\|_{1} \|\widehat{\boldsymbol{\Lambda}}_{H}\|_{1} = O_{P}(d(\log d/n)^{1/2} \cdot n^{\bar{\tau}_{V}^{*+} + \bar{\tau}_{G}^{*+}}), \\ \|(\boldsymbol{w}^{\mathsf{T}} \boldsymbol{w} - \boldsymbol{\Sigma}_{w}) \widehat{\boldsymbol{\Lambda}}_{H}\|_{\mathsf{MAX}} &\leq \|\boldsymbol{w}^{\mathsf{T}} \boldsymbol{w} - \boldsymbol{\Sigma}_{w}\|_{\mathsf{MAX}} \|\widehat{\boldsymbol{\Lambda}}_{H}\|_{1} = O_{P}(d(\log d/n)^{1/2} \cdot n^{2\bar{\tau}_{V}^{*+}}), \end{split}$$

and

$$\|\boldsymbol{\Sigma}_{w}\widehat{\boldsymbol{\Lambda}}_{H}\|_{\mathsf{MAX}} \leq \|\boldsymbol{\Sigma}_{w}\|_{\infty}\|\widehat{\boldsymbol{\Lambda}}_{H}\|_{\mathsf{MAX}} = O_{P}(d^{1/2}m_{w,nd}),$$

since $\|\mathbf{\Lambda}_H\|_{\mathsf{MAX}} = O(1)$, $\|\mathbf{\hat{\pi}}^{\mathsf{T}}\mathbf{w}\|_1 \leq K_H \|\mathbf{\hat{\pi}}^{\mathsf{T}}\mathbf{w}\|_{\mathsf{MAX}} = O_P\left((\log d/n)^{1/2} \cdot n^{\overline{\tau}_V^{*+} + \overline{\tau}_G^{*+}}\right)$ by Lemma B.4, $\|\mathbf{\hat{\Lambda}}_H\|_1 \leq d^{1/2}\|\mathbf{\hat{\Lambda}}_H\|_F = dK_H$, and $\|\mathbf{\hat{\Lambda}}_H\|_{\mathsf{MAX}} \leq \|\mathbf{\hat{\Lambda}}_H\|_F = d^{1/2}K_H$. Therefore

$$\left\| \left(\widehat{\mathbf{\Lambda}}_H - \mathbf{\Lambda}_H \mathbf{D}_H \mathbf{R} \right) \widehat{\mathbf{D}}_{x, K_H} \right\|_{\mathsf{MAX}} = O_P \left(d(\log d/n)^{1/2} \cdot n^{\overline{\tau}_V^{*+} + \overline{\tau}_G^{*+} \vee \overline{\tau}_V^{*+}} + d^{1/2} m_{w, nd} \right) = O_P (d \cdot a_{nd}).$$
(A.15)

Since $\|\widehat{\mathbf{D}}_{x,K_H}^{-1}\|_{\mathsf{sp}} = O_P(d^{-1}n^{-2\underline{\tau}_G^{*-}})$ by Lemma A.1(ii), we can prove the result by noting that

$$n^{-2\underline{\tau}_G^{*-}} \cdot a_{nd} = o_P(1)$$

when $n^{1+4\underline{\tau}_{G}^{*-}-2\bar{\tau}_{V}^{*+}-2(\bar{\tau}_{G}^{*+}\vee\bar{\tau}_{V}^{*+})}/(\log d) \to \infty$ and $m_{w,nd}/(d^{1/2}n^{2\underline{\tau}_{G}^{*-}}) \to 0.$

For part (ii), noting that

$$d^{-1}\widehat{\mathbf{\Lambda}}_{H}^{\mathsf{T}}\widehat{\mathbf{\Lambda}}_{H} = \mathbf{I}_{K_{H}} \text{ and } \left\|\mathbf{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R} - \widehat{\mathbf{\Lambda}}_{H}\right\|_{\mathsf{sp}} \leq (dK_{H})^{1/2} \left\|\mathbf{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R} - \widehat{\mathbf{\Lambda}}_{H}\right\|_{\mathsf{MAX}},$$

we have

$$\|d^{-1}\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R} - \mathbf{I}_{K_{H}}\|_{\mathsf{sp}} = d^{-1}\|\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R} - \widehat{\boldsymbol{\Lambda}}_{H}^{\mathsf{T}}\widehat{\boldsymbol{\Lambda}}_{H}\|_{\mathsf{sp}}$$

$$\leq 2d^{-1}\|\widehat{\boldsymbol{\Lambda}}_{H}\|_{\mathsf{sp}}\|\boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R} - \widehat{\boldsymbol{\Lambda}}_{H}\|_{\mathsf{sp}} + d^{-1}\|\boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R} - \widehat{\boldsymbol{\Lambda}}_{H}\|_{\mathsf{sp}}^{2}$$

$$= O_{P}(n^{-2\underline{\tau}_{G}^{*-}} \cdot a_{nd}) = o_{P}(1). \qquad (A.16)$$

Then, by triangle inequality, we have $|||(d^{-1}\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R})||_{\mathsf{sp}} - 1| = o_{P}(1)$. Since

$$\|d^{-1}\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R}\|_{\mathsf{sp}}\|(\mathbf{D}_{H}\mathbf{R}\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H})^{-1}\|_{\mathsf{sp}} \ge \|\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\boldsymbol{\Lambda}_{H}/d\|_{\mathsf{sp}},\tag{A.17}$$

by Assumption 4, we have $\|(\mathbf{D}_H \mathbf{R} \mathbf{R}^{\mathsf{T}} \mathbf{D}_H)^{-1}\|_{\mathsf{sp}} = \|(\mathbf{D}_H \mathbf{R})^{-1}\|_{\mathsf{sp}}^2 = O_P(1).$ One the other hand, by (A.2) and (A.16),

$$\|(\mathbf{D}_H\mathbf{R})^{-1}(\frac{\mathbf{\Lambda}_H^{\mathsf{T}}\mathbf{\Lambda}_H}{d})^{-1}(\mathbf{R}^{\mathsf{T}}\mathbf{D}_H)^{-1} - \mathbf{I}_{K_H}\|_{\mathsf{sp}} \le \frac{\|d^{-1}\mathbf{R}^{\mathsf{T}}\mathbf{D}_H\mathbf{\Lambda}_H^{\mathsf{T}}\mathbf{\Lambda}_H\mathbf{D}_H\mathbf{R} - \mathbf{I}_{K_H}\|_{\mathsf{sp}}}{1 - \|d^{-1}\mathbf{R}^{\mathsf{T}}\mathbf{D}_H\mathbf{\Lambda}_H^{\mathsf{T}}\mathbf{\Lambda}_H\mathbf{D}_H\mathbf{R} - \mathbf{I}_{K_H}\|_{\mathsf{sp}}} = o_P(1).$$

Then following the same argument as in (A.17), we can prove $\|\mathbf{D}_H \mathbf{R}\|_{sp} = O_P(1)$.

As for (A.12), since $\mathbf{R}^{\mathsf{T}}(\hbar^{\mathsf{T}}\hbar)^{-1}\mathbf{R}\widehat{\mathbf{D}}_{x,K_{H}} = \mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\mathbf{\Lambda}_{H}^{\mathsf{T}}\widehat{\mathbf{\Lambda}}_{H}$, by (A.10) and (A.16), we have

$$\begin{split} \|d^{-1}\widehat{\mathbf{D}}_{x,K_{H}}^{1/2}\mathbf{R}^{\mathsf{T}}(\boldsymbol{\mathbb{A}}^{\mathsf{T}}\boldsymbol{\mathbb{A}})^{-1}\mathbf{R}\widehat{\mathbf{D}}_{x,K_{H}}^{1/2} - \mathbf{I}_{K_{H}}\|_{\mathsf{sp}} \\ &= \|d^{-1}\mathbf{R}^{\mathsf{T}}(\boldsymbol{\mathbb{A}}^{\mathsf{T}}\boldsymbol{\mathbb{A}})^{-1}\mathbf{R}\widehat{\mathbf{D}}_{x,K_{H}} - \mathbf{I}_{K_{H}}\|_{\mathsf{sp}} \\ &\leq \|d^{-1}\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R} - \mathbf{I}_{K_{H}}\|_{\mathsf{sp}} + d^{-1}\|\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\|_{\mathsf{sp}}\|\widehat{\boldsymbol{\Lambda}}_{H} - \boldsymbol{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R}\|_{\mathsf{sp}} \\ &= O_{P}(n^{-2\underline{\tau}_{G}^{*-}} \cdot a_{nd}) = o_{P}(1). \end{split}$$

Then by Lemmas B.1 and B.2, we can prove $\| \mathbf{h}^{\mathsf{T}} \mathbf{h} \|_{\mathsf{sp}} = O_P(1)$, and therefore we have (A.12).

For part (iii), we use the following decomposition

$$\widehat{\boldsymbol{\mathcal{H}}}^{\mathsf{T}} - \mathbf{R}^{-1} \boldsymbol{\mathcal{H}}^{\mathsf{T}} = d^{-1} \widehat{\boldsymbol{\Lambda}}_{H}^{\mathsf{T}} \left(\boldsymbol{\Lambda}_{H} \mathbf{D}_{H} \mathbf{R} - \widehat{\boldsymbol{\Lambda}}_{H} \right) \mathbf{R}^{-1} \boldsymbol{\mathcal{H}}^{\mathsf{T}} - d^{-1} \left(\boldsymbol{\Lambda}_{H} \mathbf{D}_{H} \mathbf{R} - \widehat{\boldsymbol{\Lambda}}_{H} \right)^{\mathsf{T}} \boldsymbol{w}^{\mathsf{T}} + d^{-1} \mathbf{R}^{\mathsf{T}} \mathbf{D}_{H} \boldsymbol{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{w}^{\mathsf{T}}.$$
(A.18)

For the first term on the right hand side (RHS) of (A.18), we have

$$\begin{split} \|\widehat{\mathbf{\Lambda}}_{H}^{\mathsf{T}}\left(\mathbf{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R}-\widehat{\mathbf{\Lambda}}_{H}\right)\mathbf{R}^{-1}\mathscr{K}^{\mathsf{T}}\|_{\mathsf{sp}} \\ &\leq \|\widehat{\mathbf{\Lambda}}_{H}^{\mathsf{T}}\|\left\|\left(\mathbf{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R}-\widehat{\mathbf{\Lambda}}_{H}\right)\widehat{\mathbf{D}}_{x,K_{H}}\right\|_{\mathsf{sp}}\left\|\widehat{\mathbf{D}}_{x,K_{H}}^{-1/2}\right\|_{\mathsf{sp}}\left\|(\mathbf{R}\widehat{\mathbf{D}}_{x,K_{H}}^{1/2})^{-1}\right\|_{\mathsf{sp}}\|\mathscr{K}^{\mathsf{T}}\|_{\mathsf{sp}} \\ &= O_{P}(d^{1/2})\cdot O_{P}(d\cdot a_{nd})\cdot O_{P}(d^{-1/2}n^{-\mathscr{I}_{G}^{*-}})\cdot O_{P}(1)\cdot O_{P}(1) \\ &= O_{P}(dn^{-\mathscr{I}_{G}^{*-}}a_{nd}). \end{split}$$
(A.19)

For the second term on the RHS of (A.18), when $n^{1+4\underline{\tau}_G^{*-}-4\overline{\tau}_V^{*+}}/\log d \to \infty$, we have

$$\begin{aligned} \left\| \left(\mathbf{\Lambda}_{H} \mathbf{D}_{H} \mathbf{R} - \widehat{\mathbf{\Lambda}}_{H} \right)^{\mathsf{T}} \boldsymbol{w} \right\|_{\mathsf{sp}} &\leq \left\| \mathbf{\Lambda}_{H} \mathbf{D}_{H} \mathbf{R} - \widehat{\mathbf{\Lambda}}_{H} \right\|_{\mathsf{sp}} \| \boldsymbol{w} \|_{\mathsf{sp}} \\ &= O_{P}(d^{1/2} n^{-2\underline{\tau}_{G}^{*-}} a_{nd}) \cdot O_{P}(d^{1/2} (\log d/n)^{1/4} \cdot n^{\overline{\tau}_{V}^{*+}} + m_{w,nd}^{1/2}) \\ &= O_{P}(dn^{-\underline{\tau}_{G}^{*-}} a_{nd}) \cdot O_{P}(n^{-\underline{\tau}_{G}^{*-}} (\log d/n)^{1/4} \cdot n^{\overline{\tau}_{V}^{*+}} + m_{w,nd}^{1/2} d^{-1/2}) \\ &= o_{P}(dn^{-\underline{\tau}_{G}^{*-}} a_{nd}). \end{aligned}$$
(A.20)

For the last term on the RHS of (A.18), by Lemma B.4(v), we have

$$d^{-1} \| \mathbf{R}^{\mathsf{T}} \mathbf{D}_{H} \mathbf{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{w}^{\mathsf{T}} \|_{\mathsf{sp}} \leq d^{-1} \| \mathbf{R}^{\mathsf{T}} \mathbf{D}_{H} \|_{\mathsf{sp}} \| \mathbf{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{w}^{\mathsf{T}} \|_{\mathsf{sp}}$$

$$= d^{-1} \cdot O_{P}(1) \cdot O_{P} \left(d^{1/2} (\log d/n)^{1/4} \cdot n^{\bar{\tau}_{V}^{*+}} m_{w,nd}^{1/2} + d^{1/2} m_{w,nd}^{1/2} \right)$$

$$= o_{P}(n^{-\underline{\tau}_{G}^{*-}} a_{nd}), \qquad (A.21)$$

when $d^{-1/2}(\log d/n)^{1/4} \cdot n^{\bar{\tau}_V^{*+}} m_{w,nd}^{1/2} = o((\log d/n)^{1/2} \cdot n^{\bar{\tau}_V^{*+} + \bar{\tau}_G^{*+} \vee \bar{\tau}_V^{*-} - \underline{\tau}_G^{*-}})$, or equivalently, $n^{1-4(\bar{\tau}_G^{*+} \vee \bar{\tau}_V^{*+}) + 4\underline{\tau}_G^{*-}} = o(d^2 \log d/m_{w,nd}^2)$. Combing (A.19)–(A.21), we have $\|\hat{\mathscr{R}}^{\mathsf{T}} - \mathbf{R}^{-1}\mathscr{R}^{\mathsf{T}}\|_{\mathsf{sp}} = O_P(n^{-\underline{\tau}_G^{*-}} a_{nd})$, which completes the proof of Lemma A.2.

PROOF OF THEOREM 3.1. (i) Note that $\hat{\hat{\kappa}}^* = \hat{\hat{\kappa}}(\hat{\hat{\kappa}}^{\mathsf{T}}\hat{\hat{\kappa}})^{-1/2} = d^{1/2}\hat{\hat{\kappa}}\hat{\mathbf{D}}_{x,K_H}^{-1/2}$ and

$$(\mathbf{R}^*)^{-1} = d^{1/2} \widehat{\mathbf{D}}_{x,K_H}^{-1/2} \mathbf{R}^{-1} (\mathscr{K}^{\mathsf{T}} \mathscr{R})^{1/2}.$$
 (A.22)

By Lemma A.2(iii), we have

$$\begin{split} \left\| \widehat{\boldsymbol{\lambda}}^{*^{\mathsf{T}}} - (\mathbf{R}^{*})^{-1} (\boldsymbol{\lambda}^{\mathsf{T}} \boldsymbol{\lambda})^{-1/2} \boldsymbol{\lambda}^{\mathsf{T}} \right\|_{\mathsf{MAX}} &= \left\| \widehat{\boldsymbol{\lambda}}^{*^{\mathsf{T}}} - d^{1/2} \widehat{\mathbf{D}}_{\boldsymbol{x},\boldsymbol{K}_{H}}^{-1/2} \mathbf{R}^{-1} (\boldsymbol{\lambda}^{\mathsf{T}} \boldsymbol{\lambda})^{1/2} (\boldsymbol{\lambda}^{\mathsf{T}} \boldsymbol{\lambda})^{-1/2} \boldsymbol{\lambda}^{\mathsf{T}} \right\|_{\mathsf{MAX}} \\ &= \left\| \widehat{\boldsymbol{\lambda}}^{*^{\mathsf{T}}} - d^{1/2} \widehat{\mathbf{D}}_{\boldsymbol{x},\boldsymbol{K}_{H}}^{-1/2} \mathbf{R}^{-1} \boldsymbol{\lambda}^{\mathsf{T}} \right\|_{\mathsf{MAX}} \\ &\leq d^{1/2} \| \widehat{\mathbf{D}}_{\boldsymbol{x},\boldsymbol{K}_{H}}^{-1/2} \|_{\mathsf{sp}} \| \widehat{\boldsymbol{\lambda}}^{\mathsf{T}} - \mathbf{R}^{-1} \boldsymbol{\lambda}^{\mathsf{T}} \|_{\mathsf{sp}} \\ &= O_{P} \left(n^{-2\underline{\tau}_{G}^{*-}} \cdot a_{nd} \right). \end{split}$$

(ii) Following (A.15) and noting that $\widehat{\Lambda}_{H}^{*} = \widehat{\Lambda}_{H} (\widehat{\lambda}^{\mathsf{T}} \widehat{\lambda})^{1/2} = d^{-1/2} \widehat{\Lambda}_{H} \widehat{\mathbf{D}}_{x,K_{H}}^{1/2}$, we have

$$\left\|\widehat{\mathbf{\Lambda}}_{H}^{*} - d^{-1/2}\mathbf{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R}\widehat{\mathbf{D}}_{x,K_{H}}^{1/2}\right\|_{\mathsf{MAX}} = O_{P}\left(n^{-2\underline{\tau}_{G}^{*-}} \cdot a_{nd}\right).$$

Using the notation of \mathbf{R}^* , it can be equivalently written as

$$\left\|\widehat{\boldsymbol{\Lambda}}_{H}^{*}-\boldsymbol{\Lambda}_{H}\mathbf{D}_{H}(\boldsymbol{\varkappa}^{\mathsf{T}}\boldsymbol{\varkappa})^{1/2}\mathbf{R}^{*}\right\|_{\mathsf{MAX}}=O_{P}\left(n^{-2\underline{\tau}_{G}^{*-}}\cdot a_{nd}\right).$$

(iii) For the first part, it is obvious that $\hat{k}^* \hat{\Lambda}_H^{*^{\mathsf{T}}} = \hat{k} \hat{\Lambda}_H^{\mathsf{T}}$. For the second part, we have

$$\begin{split} & \left\| \widehat{\mathcal{R}}^* \widehat{\boldsymbol{\Lambda}}_H^{*^{\mathsf{T}}} - \mathscr{R} \mathbf{D}_H \mathbf{\Lambda}_H^{\mathsf{T}} \right\|_{\mathsf{MAX}} \\ & \leq \left\| \widehat{\mathcal{R}}^{*^{\mathsf{T}}} - (\mathbf{R}^*)^{-1} (\mathscr{R}^{\mathsf{T}} \mathscr{R})^{-1/2} \mathscr{R}^{\mathsf{T}} \right\|_1 \left\| \mathbf{\Lambda}_H \right\|_{\mathsf{MAX}} \left\| \mathbf{D}_H (\mathscr{R}^{\mathsf{T}} \mathscr{R})^{1/2} \mathbf{R}^* \right\|_1 \\ & + \left\| (\mathbf{R}^*)^{-1} (\mathscr{R}^{\mathsf{T}} \mathscr{R})^{-1/2} \mathscr{R}^{\mathsf{T}} \right\|_1 \left\| \widehat{\boldsymbol{\Lambda}}_H^* - \mathbf{\Lambda}_H \mathbf{D}_H (\mathscr{R}^{\mathsf{T}} \mathscr{R})^{1/2} \mathbf{R}^* \right\|_{\mathsf{MAX}} \\ & + \left\| \widehat{\mathscr{R}}^{*^{\mathsf{T}}} - (\mathbf{R}^*)^{-1} (\mathscr{R}^{\mathsf{T}} \mathscr{R})^{-1/2} \mathscr{R}^{\mathsf{T}} \right\|_1 \left\| \widehat{\boldsymbol{\Lambda}}_H^* - \mathbf{\Lambda}_H \mathbf{D}_H (\mathscr{R}^{\mathsf{T}} \mathscr{R})^{1/2} \mathbf{R}^* \right\|_{\mathsf{MAX}} \\ & = O_P \left(n^{\overline{\tau}_G^{*+} - 2\underline{\tau}_G^{*-}} \cdot a_{nd} \right), \end{split}$$

as $\|\mathbf{\Lambda}_H\|_{\mathsf{MAX}} = O(1), \|\mathbf{D}_H(\mathscr{K}^{\mathsf{T}}\mathscr{K})^{1/2}\mathbf{R}^*\|_1 = O_P(n^{\overline{\tau}_G^{*+}}), \text{ and}$

$$\left\|\widehat{\mathscr{R}}^{*^{\mathsf{T}}} - (\mathbf{R}^{*})^{-1}(\mathscr{R}^{\mathsf{T}}\mathscr{R})^{-1/2}\mathscr{R}^{\mathsf{T}}\right\|_{1} \leq K_{H}^{1/2} \left\|\widehat{\mathscr{R}}^{*^{\mathsf{T}}} - (\mathbf{R}^{*})^{-1}(\mathscr{R}^{\mathsf{T}}\mathscr{R})^{-1/2}\mathscr{R}^{\mathsf{T}}\right\|_{\mathsf{sp}} = O_{P}(n^{\overline{\tau}_{G}^{*+}} \cdot a_{nd}).$$

Thus, we obtain the uniformly convergence rate for the common components.

(iv) We next prove that \mathbf{R}^* is an asymptotically orthogonal matrix. By (A.22), we have

$$\mathbf{R}^{*^{\mathsf{T}}}\mathbf{R}^{*} = d^{-1}\widehat{\mathbf{D}}_{x,K_{H}}^{1/2}\mathbf{R}^{\mathsf{T}}(\mathscr{R}^{\mathsf{T}}\mathscr{R})^{-1}\mathbf{R}\widehat{\mathbf{D}}_{x,K_{H}}^{1/2} = d^{-1}\widehat{\mathbf{D}}_{x,K_{H}}^{1/2}\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\boldsymbol{\Lambda}_{H}^{\mathsf{T}}\widehat{\boldsymbol{\Lambda}}_{H}\widehat{\mathbf{D}}_{x,K_{H}}^{-1/2}.$$

Thus by (A.10) and (A.16), we have

$$\begin{aligned} \|\mathbf{R}^{*^{\mathsf{T}}}\mathbf{R}^{*} - \mathbf{I}_{K_{H}}\|_{\mathsf{sp}} \\ &\leq \|d^{-1}\widehat{\mathbf{D}}_{x,K_{H}}^{1/2}\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\mathbf{\Lambda}_{H}^{\mathsf{T}}\widehat{\mathbf{\Lambda}}_{H}\widehat{\mathbf{D}}_{x,K_{H}}^{-1/2} - d^{-1}\widehat{\mathbf{D}}_{x,K_{H}}^{1/2}\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\mathbf{\Lambda}_{H}^{\mathsf{T}}\mathbf{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R}\widehat{\mathbf{D}}_{x,K_{H}}^{-1/2}\|_{\mathsf{sp}} \\ &+ \|d^{-1}\widehat{\mathbf{D}}_{x,K_{H}}^{1/2}\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\mathbf{\Lambda}_{H}^{\mathsf{T}}\mathbf{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R}\widehat{\mathbf{D}}_{x,K_{H}}^{-1/2} - \mathbf{I}_{K_{H}}\|_{\mathsf{sp}} \\ &= \|d^{-1}\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\mathbf{\Lambda}_{H}^{\mathsf{T}}\widehat{\mathbf{\Lambda}}_{H} - d^{-1}\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\mathbf{\Lambda}_{H}^{\mathsf{T}}\mathbf{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R}\|_{\mathsf{sp}} \\ &+ \|d^{-1}\mathbf{R}^{\mathsf{T}}\mathbf{D}_{H}\mathbf{\Lambda}_{H}^{\mathsf{T}}\mathbf{\Lambda}_{H}\mathbf{D}_{H}\mathbf{R} - \mathbf{I}_{K_{H}}\|_{\mathsf{sp}} \\ &= O_{P}(n^{-2\underline{\tau}_{G}^{*-}}a_{nd}). \end{aligned}$$
(A.23)

We thus complete the proof of Theorem 3.1.

To prove Lemma 3.1, we need some intermediate estimators or infeasible estimators related to $\hat{\boldsymbol{\beta}}$ and $\hat{\boldsymbol{\beta}}_{\perp}$. Recall that $\hat{\boldsymbol{\beta}}_{\perp}$ is the matrix of eigenvectors associated with the largest K_F eigenvalues of $\hat{\mathbf{S}}_{HH} := n^{-1} \widehat{\mathscr{H}}^{*c^{\mathsf{T}}} \widehat{\mathscr{H}}^{*c}$, and that $\hat{\boldsymbol{\beta}}$ is the matrix of eigenvectors associated with the rest of the K_G eigenvalues. For a $K_H \times K_H$ matrix, $\boldsymbol{\Xi}$, we define $\mathbf{S}_{HH}^{\boldsymbol{\Xi}} := n^{-1} \boldsymbol{\Xi} \mathscr{H}^{c^{\mathsf{T}}} \mathscr{H}^{c} \boldsymbol{\Xi}^{\mathsf{T}}$, where $\mathscr{H}^c = \mathscr{H} - \overline{\mathscr{H}}$ and $\overline{\mathscr{H}} = n^{-1} \mathbf{1}_n \sum_{s=1}^n \mathbf{H}_{s\Delta}^{\mathsf{T}}$. Replacing $\hat{\mathbf{S}}_{HH}$ with $\mathbf{S}_{HH}^{\boldsymbol{\Xi}}$, we can obtain the infeasible estimators, $\boldsymbol{\beta}_{\perp}^{\boldsymbol{\Xi}}$ and $\boldsymbol{\beta}^{\boldsymbol{\Xi}}$. For simplicity, we use \mathbf{S}_{HH} to denote $\mathbf{S}_{HH}^{\mathbf{I}_{KH}}$. Later on, we will determine a proper choice of $\boldsymbol{\Xi}$.

Lemma A.3. Suppose that Assumptions 1, 3 and 6 are satisfied. If the eigenvalues of $\Xi\Xi^{\mathsf{T}}$ are bounded away from zero and infinity uniformly with probability approaching one, then β_{\perp}^{Ξ} and β^{Ξ} are super-consistent in the sense that

$$\boldsymbol{\beta}^{\boldsymbol{\Xi}} - \boldsymbol{\Xi}^{\mathsf{T}} \boldsymbol{\beta} [\boldsymbol{\beta}^{\mathsf{T}} \boldsymbol{\Xi} \boldsymbol{\Xi}^{\mathsf{T}} \boldsymbol{\beta}]^{-1} \boldsymbol{\beta}^{\mathsf{T}} \boldsymbol{\Xi} \boldsymbol{\beta}^{\boldsymbol{\Xi}} = O_P(n^{-1}), \qquad (A.24)$$

and

$$\boldsymbol{\beta}_{\perp}^{\boldsymbol{\Xi}} - \boldsymbol{\Xi}^{-1} \boldsymbol{\beta}_{\perp} [\boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\boldsymbol{\Xi}^{\mathsf{T}})^{-1} \boldsymbol{\Xi}^{-1} \boldsymbol{\beta}_{\perp}]^{-1} \boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\boldsymbol{\Xi}^{\mathsf{T}})^{-1} \boldsymbol{\beta}_{\perp}^{\boldsymbol{\Xi}} = O_P(n^{-1}).$$
(A.25)

Proof. We decompose β^{Ξ} in the directions of $\Xi^{\mathsf{T}}\beta$ and $\Xi^{-1}\beta_{\perp}$ (which are orthogonal) as

$$\boldsymbol{\beta}^{\Xi} = \Xi^{\mathsf{T}} \boldsymbol{\beta} [\boldsymbol{\beta}^{\mathsf{T}} \Xi \Xi^{\mathsf{T}} \boldsymbol{\beta}]^{-1} \boldsymbol{\beta}^{\mathsf{T}} \Xi \boldsymbol{\beta}^{\Xi} + \Xi^{-1} \boldsymbol{\beta}_{\perp} [\boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\Xi^{\mathsf{T}})^{-1} \Xi^{-1} \boldsymbol{\beta}_{\perp}]^{-1} \boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\Xi^{\mathsf{T}})^{-1} \boldsymbol{\beta}^{\Xi}.$$
(A.26)

Note that $\boldsymbol{\beta}^{\Xi}$ satisfies $\boldsymbol{\Xi}^{-1}\mathbf{S}_{HH}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\boldsymbol{\beta}^{\Xi} = \boldsymbol{\beta}^{\Xi}\mathbf{D}_{S}^{\Xi}$, where \mathbf{D}_{S}^{Ξ} is a $K_{G} \times K_{G}$ diagonal matrix with the diagonal elements being the K_{G} smallest eigenvalues of $\boldsymbol{\Xi}^{-1}\mathbf{S}_{HH}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}$ arranged in a descending order. Using (A.26) and the equality $\mathbf{I}_{K_{H}} = \boldsymbol{\beta}_{\perp}\boldsymbol{\beta}_{\perp}^{\mathsf{T}} + \boldsymbol{\beta}\boldsymbol{\beta}^{\mathsf{T}}$, we have

$$\begin{split} \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \Xi^{-1} \boldsymbol{\beta}^{\Xi} \mathbf{D}_{S}^{\Xi} = & \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \Xi\Xi^{-1} \mathbf{S}_{HH} (\Xi^{\mathsf{T}})^{-1} \boldsymbol{\beta}^{\Xi} \\ = & \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta} [\boldsymbol{\beta}^{\mathsf{T}} \Xi\Xi^{\mathsf{T}} \boldsymbol{\beta}]^{-1} \boldsymbol{\beta}^{\mathsf{T}} \Xi\boldsymbol{\beta}^{\Xi} \\ & + & \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \mathbf{I}_{K_{H}} (\Xi^{\mathsf{T}})^{-1} \Xi^{-1} \boldsymbol{\beta}_{\perp} [\boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\Xi^{\mathsf{T}})^{-1} \Xi^{-1} \boldsymbol{\beta}_{\perp}]^{-1} \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \Xi^{-1} \boldsymbol{\beta}^{\Xi} \\ = & \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta} [\boldsymbol{\beta}^{\mathsf{T}} \Xi\Xi^{\mathsf{T}} \boldsymbol{\beta}]^{-1} \boldsymbol{\beta}^{\mathsf{T}} \Xi\boldsymbol{\beta}^{\Xi} \\ & + & \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta} \boldsymbol{\beta}^{\mathsf{T}} (\Xi^{\mathsf{T}})^{-1} \Xi^{-1} \boldsymbol{\beta}_{\perp} [\boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\Xi^{\mathsf{T}})^{-1} \Xi^{-1} \boldsymbol{\beta}_{\perp}]^{-1} \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \Xi^{-1} \boldsymbol{\beta}^{\Xi} \\ & + & \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta}_{\perp} \boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\Xi^{\mathsf{T}})^{-1} \Xi^{-1} \boldsymbol{\beta}_{\perp} [\boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\Xi^{\mathsf{T}})^{-1} \Xi^{-1} \boldsymbol{\beta}_{\perp}]^{-1} \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \Xi^{-1} \boldsymbol{\beta}^{\Xi} \\ & = & \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta} [\boldsymbol{\beta}^{\mathsf{T}} \Xi\Xi^{\mathsf{T}} \boldsymbol{\beta}]^{-1} \boldsymbol{\beta}^{\mathsf{T}} \Xi\boldsymbol{\beta}^{\Xi} \\ & + & \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta} \boldsymbol{\beta}^{\mathsf{T}} (\Xi^{\mathsf{T}})^{-1} \Xi^{-1} \boldsymbol{\beta}_{\perp} [\boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\Xi^{\mathsf{T}})^{-1} \Xi^{-1} \boldsymbol{\beta}_{\perp}]^{-1} \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \Xi^{-1} \boldsymbol{\beta}^{\Xi} \\ & + & \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta} \boldsymbol{\beta}^{\mathsf{T}} \Xi^{\mathsf{T}} \boldsymbol{\beta}]^{-1} \boldsymbol{\beta}^{\mathsf{T}} \Xi\boldsymbol{\beta}^{\Xi} \\ & + & \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta} \boldsymbol{\beta}^{\mathsf{T}} (\Xi^{\mathsf{T}})^{-1} \Xi^{-1} \boldsymbol{\beta}_{\perp} [\boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\Xi^{\mathsf{T}})^{-1} \Xi^{-1} \boldsymbol{\beta}_{\perp}]^{-1} \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \Xi^{-1} \boldsymbol{\beta}^{\Xi} \end{split}$$

Vectorizing this expression, we have

$$\operatorname{vec}(\boldsymbol{\beta}_{\perp}^{\mathsf{T}}\boldsymbol{\Xi}^{-1}\boldsymbol{\beta}^{\boldsymbol{\Xi}}) = \left\{ \mathbf{D}_{S}^{\boldsymbol{\Xi}} \otimes \mathbf{I}_{K_{F}} - \mathbf{I}_{K_{G}} \otimes \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta}_{\perp} - \mathbf{I}_{K_{G}} \otimes \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta} \boldsymbol{\beta}^{\mathsf{T}}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1} \boldsymbol{\Xi}^{-1} \boldsymbol{\beta}_{\perp} [\boldsymbol{\beta}_{\perp}^{\mathsf{T}}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1} \boldsymbol{\Xi}^{-1} \boldsymbol{\beta}_{\perp}]^{-1} \right\}^{-1} \\ \cdot \operatorname{vec}\left(\boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta} [\boldsymbol{\beta}^{\mathsf{T}} \boldsymbol{\Xi} \boldsymbol{\Xi}^{\mathsf{T}} \boldsymbol{\beta}]^{-1} \boldsymbol{\beta}^{\mathsf{T}} \boldsymbol{\Xi} \boldsymbol{\beta}^{\boldsymbol{\Xi}}\right).$$
(A.27)

Recall that $\boldsymbol{\beta} = (\mathbf{O}_{K_G \times K_F} \quad \mathbf{I}_{K_G})^{\mathsf{T}}$ and $\boldsymbol{\beta}_{\perp} = (\mathbf{I}_{K_F} \quad \mathbf{O}_{K_F \times K_G})^{\mathsf{T}}$. By Lemma B.5, we have

$$\boldsymbol{\beta}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta} = n^{-2} \sum_{s=1}^{n} \boldsymbol{G}_{s\Delta}^{c} \boldsymbol{G}_{s\Delta}^{c^{\mathsf{T}}} = O_{P}(n^{-1}),$$
$$\boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta}_{\perp} = n^{-1} \sum_{s=1}^{n} \boldsymbol{F}_{s\Delta}^{c} \boldsymbol{F}_{s\Delta}^{c^{\mathsf{T}}} \text{ is bounded away from zero.}$$
$$\boldsymbol{\beta}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta}_{\perp} = n^{-3/2} \sum_{s=1}^{n} \boldsymbol{G}_{s\Delta}^{c} \boldsymbol{F}_{s\Delta}^{c^{\mathsf{T}}} = O_{P}(n^{-1}),$$

where $\mathbf{G}_{s\Delta}^{c} = \mathbf{G}_{s\Delta} - n^{-1} \sum_{s=1}^{n} \mathbf{G}_{s\Delta}$ and $\mathbf{F}_{s\Delta}^{c} = \mathbf{F}_{s\Delta} - n^{-1} \sum_{s=1}^{n} \mathbf{F}_{s\Delta}$. Thus, only the first block, $\boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta}_{\perp}$, of the matrix $\mathbf{S}_{HH} = \begin{bmatrix} \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta}_{\perp} & \boldsymbol{\beta}_{\perp}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta} \\ \boldsymbol{\beta}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta}_{\perp} & \boldsymbol{\beta}^{\mathsf{T}} \mathbf{S}_{HH} \boldsymbol{\beta} \end{bmatrix}$ does not converge to zero. Therefore, $\mathbf{D}_{S}^{\Xi} = o_{P}(1)$, and we have $\boldsymbol{\beta}_{\perp}^{\mathsf{T}} \Xi^{-1} \boldsymbol{\beta}^{\Xi} = O_{P}(n^{-1})$. Then using (A.26) again, we can prove the consistency of $\boldsymbol{\beta}^{\Xi}$. Using the same argument, we can prove the consistency of $\boldsymbol{\beta}_{\perp}^{\Xi}$. When $\boldsymbol{\Xi} = \mathbf{I}_{K_H}$, the results degenerate to Lemma 1 of Harris (1997). When $\boldsymbol{\Xi} = (\boldsymbol{\mathbb{A}}^{\mathsf{T}}\boldsymbol{\mathbb{A}})^{1/2}(\mathbf{R}^*)$, and replacing \mathbf{S}_{HH} with $\boldsymbol{\Xi}\widehat{\mathbf{S}}_{HH}\boldsymbol{\Xi}^{\mathsf{T}}$, we can prove Lemma 3.1.

PROOF OF LEMMA 3.1. Note that $\hat{\boldsymbol{\beta}}$ satisfies $\hat{\mathbf{S}}_{HH}\hat{\boldsymbol{\beta}} = \hat{\boldsymbol{\beta}}\hat{\mathbf{D}}_S$, where $\hat{\mathbf{D}}_S$ is a $K_G \times K_G$ diagonal matrix with the diagonal elements being the K_G smallest eigenvalues of $\hat{\mathbf{S}}_{HH}$ arranged in a descending order. Let $\boldsymbol{\Xi} = (\boldsymbol{\mathbb{A}}^{\mathsf{T}}\boldsymbol{\mathbb{A}})^{1/2}(\mathbf{R}^*)^{\mathsf{T}}$. Following similar arguments in the proof of Lemma A.3, we have

$$\operatorname{vec}(\boldsymbol{\beta}_{\perp}^{\mathsf{T}}\boldsymbol{\Xi}^{-1}\widehat{\boldsymbol{\beta}}) = \left\{ \widehat{\mathbf{D}}_{S}^{\boldsymbol{\Xi}} \otimes \mathbf{I}_{K_{F}} - \mathbf{I}_{K_{G}} \otimes \boldsymbol{\beta}_{\perp}^{\mathsf{T}} [\boldsymbol{\Xi}\widehat{\mathbf{S}}_{HH}\boldsymbol{\Xi}^{\mathsf{T}}] \boldsymbol{\beta}_{\perp} - \mathbf{I}_{K_{G}} \otimes \boldsymbol{\beta}_{\perp}^{\mathsf{T}} [\boldsymbol{\Xi}\widehat{\mathbf{S}}_{HH}\boldsymbol{\Xi}^{\mathsf{T}}] \boldsymbol{\beta} \boldsymbol{\beta}^{\mathsf{T}} (\boldsymbol{\Xi}^{\mathsf{T}})^{-1} \boldsymbol{\Xi}^{-1} \boldsymbol{\beta}_{\perp} [\boldsymbol{\beta}_{\perp}^{\mathsf{T}} (\boldsymbol{\Xi}^{\mathsf{T}})^{-1} \boldsymbol{\Xi}^{-1} \boldsymbol{\beta}_{\perp}]^{-1} \right\}^{-1} \\ \cdot \operatorname{vec}\left(\boldsymbol{\beta}_{\perp}^{\mathsf{T}} [\boldsymbol{\Xi}\widehat{\mathbf{S}}_{HH}\boldsymbol{\Xi}^{\mathsf{T}}] \boldsymbol{\beta} [\boldsymbol{\beta}^{\mathsf{T}} \boldsymbol{\Xi} \boldsymbol{\Xi}^{\mathsf{T}} \boldsymbol{\beta}]^{-1} \boldsymbol{\beta}^{\mathsf{T}} \boldsymbol{\Xi} \widehat{\boldsymbol{\beta}} \right).$$
(A.28)

Using the convergence results in Lemma B.6 and $\|\widehat{\mathbf{D}}_S\|_{sp} \leq \|\mathbf{D}_S^{\Xi}\|_{sp} + \|\widehat{\mathbf{S}}_{HH} - \Xi^{-1}\mathbf{S}_{HH}(\Xi^{\mathsf{T}})^{-1}\|_{sp} = o_P(1)$, we can prove $\beta_{\perp}^{\mathsf{T}}\Xi^{-1}\widehat{\boldsymbol{\beta}} = O_P(n^{-1})$. Then following the same arguments as in the proof of Lemma A.3, we can prove the results.

PROOF OF THEOREM 3.2. (ii) and (iii) follow directly from Theorem 3.1 and Lemma 3.1 by noting that $a_{nd} > n^{-1}$.

As for (i), by Theorem 3.1 and Lemma 3.1, we have

$$\left\|\widehat{\mathscr{K}}^*\widehat{\boldsymbol{\beta}}_{\perp} - \mathscr{K}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\boldsymbol{\Xi}^{-1}\boldsymbol{\beta}_{\perp}\mathbf{Q}_{\boldsymbol{\beta}_{\perp}}\right\|_{\mathsf{sp}} = O_P\left(n^{-2\underline{\tau}_G^{*-}} \cdot a_{nd}\right).$$
(A.29)

Using $(\boldsymbol{\beta}_{\perp}\boldsymbol{\beta}_{\perp}^{\mathsf{T}} + \boldsymbol{\beta}\boldsymbol{\beta}^{\mathsf{T}}) = \mathbf{I}_{K_{H}}$, we have

$$\begin{split} & \boldsymbol{\hslash}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\boldsymbol{\Xi}^{-1}\boldsymbol{\beta}_{\perp}\mathbf{Q}_{\boldsymbol{\beta}_{\perp}} = \boldsymbol{\hslash}(\boldsymbol{\beta}_{\perp}\boldsymbol{\beta}_{\perp}^{\mathsf{T}} + \boldsymbol{\beta}\boldsymbol{\beta}^{\mathsf{T}})(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\boldsymbol{\Xi}^{-1}\boldsymbol{\beta}_{\perp}[\boldsymbol{\beta}_{\perp}^{\mathsf{T}}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\boldsymbol{\Xi}^{-1}\boldsymbol{\beta}_{\perp}]^{-1}\boldsymbol{\beta}_{\perp}^{\mathsf{T}}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\widehat{\boldsymbol{\beta}}_{\perp} \\ & = \boldsymbol{\pounds}\boldsymbol{\beta}_{\perp}^{\mathsf{T}}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\widehat{\boldsymbol{\beta}}_{\perp} + \boldsymbol{\pounds}\boldsymbol{\beta}^{\mathsf{T}}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\boldsymbol{\Xi}^{-1}\boldsymbol{\beta}_{\perp}[\boldsymbol{\beta}_{\perp}^{\mathsf{T}}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\boldsymbol{\Xi}^{-1}\boldsymbol{\beta}_{\perp}]^{-1}\boldsymbol{\beta}_{\perp}^{\mathsf{T}}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\widehat{\boldsymbol{\beta}}_{\perp}. \end{split}$$

Therefore, we only need to prove that the second term on the RHS of the second equality above is $O_P(n^{-2\mathfrak{I}_G^*}a_{nd})$. Indeed,

$$\boldsymbol{\beta}^{\mathsf{T}}(\boldsymbol{\Xi}^{\mathsf{T}})^{-1}\boldsymbol{\Xi}^{-1}\boldsymbol{\beta}_{\perp} = \boldsymbol{\beta}^{\mathsf{T}}(\boldsymbol{\hbar}^{\mathsf{T}}\boldsymbol{\hbar})^{-1/2}(\mathbf{R}^{*\mathsf{T}}\mathbf{R}^{*})^{-1}(\boldsymbol{\hbar}^{\mathsf{T}}\boldsymbol{\hbar})^{-1/2}\boldsymbol{\beta}_{\perp}$$
$$= \boldsymbol{\beta}^{\mathsf{T}}(\boldsymbol{\hbar}^{\mathsf{T}}\boldsymbol{\hbar})^{-1}\boldsymbol{\beta}_{\perp} + O_{P}(n^{-2\underline{\tau}_{G}^{*-}}a_{nd}) = O_{P}(n^{-2\underline{\tau}_{G}^{*-}}a_{nd}),$$

where the last two equalities follow from (A.23) and the result that $\mathscr{R}^{\intercal}\mathscr{R}$ converges to a block diagonal matrix at rate $(\log d/n)^{1/2}$ using Lemmas B.1–B.3. Thus we complete the proof of (i).

As for (iv), by Theorem 3.1 and Lemma 3.1, we have

$$\left\|\widehat{\mathbf{\Lambda}}_{H}^{*}\widehat{\boldsymbol{\beta}} - \mathbf{\Lambda}_{H}\mathbf{D}_{H}\Xi\Xi^{\mathsf{T}}\boldsymbol{\beta}\mathbf{Q}_{\boldsymbol{\beta}}\right\|_{\mathsf{sp}} = O_{P}\left(n^{-2\underline{\tau}_{G}^{*-}} \cdot a_{nd}\right).$$
(A.30)

Following similar arguments to the proof of part (i), we can show that $\boldsymbol{\beta}_{\perp}^{\mathsf{T}} \boldsymbol{\Xi} \boldsymbol{\Xi}^{\mathsf{T}} \boldsymbol{\beta} = O_P(n^{-2\underline{\tau}_G^{*-}} a_{nd})$ and that

$$\left\| \mathbf{\Lambda}_H \mathbf{D}_H \mathbf{\Xi} \mathbf{\Xi}^{\mathsf{T}} \boldsymbol{\beta} \mathbf{Q}_{\boldsymbol{\beta}} - \mathbf{\Lambda}_G \mathbf{D}_G \boldsymbol{\beta}^{\mathsf{T}} \mathbf{\Xi} \widehat{\boldsymbol{\beta}} \right\|_{\mathsf{sp}} = O_P \left(n^{-2\underline{\tau}_G^{*-}} \cdot a_{nd} \right).$$
(A.31)

Combining (A.30) and (A.31), we complete the proof.

PROOF OF THEOREM 3.3.

Since $\|\mathbf{L}_n\|_{sp} = 1$, using the submultiplicative property of the spectral norm, the convergence rate of the estimators for the first-differenced factors is also the convergence rate of the estimators for the cumulated factors.

Appendix B: Auxiliary Lemmas

Recall that f_t and u_t are increments of continuous-time processes, between t and $t - \Delta$, while g_t and v_t are the first-order differences of stationary time series, G_t and V_t , for $t = 0, \Delta, \ldots, n\Delta$. Lemmas B.1–B.3 give the large deviation theory for them. Specially, Lemma B.1 is for f_t and u_t only, Lemma B.2 for g_t and v_t only, and Lemma B.3 for both the continuous-time processes and the discrete-time processes.

Lemma B.1. Under Assumption 1, we have (i) $\left\|\sum_{s=1}^{n} \boldsymbol{u}_{s\Delta} \boldsymbol{u}_{s\Delta}^{\mathsf{T}} - \boldsymbol{\Sigma}_{U}\right\|_{\mathsf{MAX}} = O_{P}((\log d/n)^{1/2});$ (ii) $\left\|\sum_{s=1}^{n} \boldsymbol{f}_{s\Delta} \boldsymbol{f}_{s\Delta}^{\mathsf{T}} - \boldsymbol{\Sigma}_{F}\right\|_{\mathsf{MAX}} = O_{P}((\log d/n)^{1/2});$ (iii) $\left\|\sum_{s=1}^{n} \boldsymbol{u}_{s\Delta} \boldsymbol{f}_{s\Delta}^{\mathsf{T}}\right\|_{\mathsf{MAX}} = O_{P}((\log d/n)^{1/2}).$ (iv) In addition, if Assumptions 4 and 5 hold, we have $\left\|d^{-1}\sum_{s=1}^{n} \boldsymbol{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{u}_{s\Delta} \boldsymbol{u}_{s\Delta}^{\mathsf{T}} \boldsymbol{\Lambda}_{H} - d^{-1} \boldsymbol{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{\Sigma}_{U} \boldsymbol{\Lambda}_{H}\right\|_{\mathsf{MAX}} = O_{P}(m_{U,d}(\log d/n)^{1/2}).$

Proof. Parts (i)–(iii) are the same as Lemma 1 in Aït-Sahalia and Xiu (2017). We only prove part (iv), as parts (i)–(iii) can be proved similarly. By Bonferroni inequality and Lemma 10 of Tao et al. (2013b), we have

$$P\left(\left\|\sum_{s=1}^{n} d^{-1} \boldsymbol{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{u}_{s\Delta} \boldsymbol{u}_{s\Delta}^{\mathsf{T}} \boldsymbol{\Lambda}_{H} - d^{-1} \boldsymbol{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{\Sigma}_{U} \boldsymbol{\Lambda}_{H}\right\|_{\mathsf{MAX}} > c\right) \leq K_{h}^{2} \cdot 4 \exp(-nc^{2}/(64C_{1}))$$

for all $0 \le c \le \mu_d^2 (d^{-1} \mathbf{\Lambda}_H^{\mathsf{T}} \mathbf{\Sigma}_U \mathbf{\Lambda}_H) \cdot n^{1/2}$, where $C_1 = 8 \| d^{-1} \mathbf{\Lambda}_H^{\mathsf{T}} \mathbf{\Sigma}_U \mathbf{\Lambda}_H \|_{\mathsf{MAX}}^2$ is obtained from Lemma 3 of Fan et al. (2012). By Assumptions 1 and 4, we have that $\mu_d^2 (d^{-1} \mathbf{\Lambda}_H^{\mathsf{T}} \mathbf{\Sigma}_U \mathbf{\Lambda}_H)$ is bounded away from zero, and $C_1 \le 8 \| d^{-1} \mathbf{\Lambda}_H^{\mathsf{T}} \mathbf{\Sigma}_U \mathbf{\Lambda}_H \|_{\mathsf{sp}}^2 = O(m_{U,d}^2)$. Then using the exponential inequality and taking $x = m_{U,d} (\log d/n)^{1/2}$, we can prove the result. \Box

Lemma B.2. Under Assumption 6, we have (i) $\|n^{-1}\sum_{s=1}^{n} \boldsymbol{v}_{s\Delta} \boldsymbol{v}_{s\Delta}^{\mathsf{T}} - \boldsymbol{\Sigma}_{v}\|_{\mathsf{MAX}} = O_{P}((\log d/n)^{1/2});$ (ii) $\|n^{-1}\sum_{s=1}^{n} \boldsymbol{g}_{s\Delta} \boldsymbol{g}_{s\Delta}^{\mathsf{T}} - \boldsymbol{\Sigma}_{g}\|_{\mathsf{MAX}} = O_{P}((\log d/n)^{1/2});$ (iii) $\|n^{-1}\sum_{s=1}^{n} \boldsymbol{v}_{s\Delta} \boldsymbol{g}_{s\Delta}^{\mathsf{T}}\|_{\mathsf{MAX}} = O_{P}((\log d/n)^{1/2}).$ (iv) In addition, if Assumptions 4 and 5 hold, we have $\|(nd)^{-1}\sum_{s=1}^{n} \boldsymbol{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{v}_{s\Delta} \boldsymbol{v}_{s\Delta}^{\mathsf{T}} \boldsymbol{\Lambda}_{H} - d^{-1} \boldsymbol{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{\Sigma}_{v} \boldsymbol{\Lambda}_{H}\|_{\mathsf{MAX}} = O_{P}(m_{v,d}(\log d/n)^{1/2}).$

Proof. Parts (i)–(iii) are the same as Lemma C.3 in Fan et al. (2013). Note that, under Assumption 4(i), the mixing coefficient of $\{v_{s\Delta}\}$ is bounded by $C'_{\alpha} \exp(-s^{\gamma_2})$, for some positive constant C'_{α} . Also note that $v_{i,s\Delta}$ still satisfies the exponential-type tail condition, since

$$\max_{1 \le i \le d} P\left(|v_{i,s\Delta}| > c\right) \le \max_{1 \le i \le d} 2P\left(|V_{i,s\Delta}| > c/2\right) \\
\le 2\exp(1 - (c/(2b_1))^{\gamma_2}) \le \exp(1 - (c/b_3)^{\gamma_4}),$$
(B.1)

for $1 \leq i \leq d, s = 1, \dots, n$, and c > 0, where $\gamma_4 \in (0, \gamma_2)$ and $b_3 > 2b_1 \max\{(\gamma_4/\gamma_2)^{1/\gamma_2}, (1+\log 2)^{1/\gamma_2}\}$, and the last inequality is shown in the proof of Lemma C.2 of Fan et al. (2011). Again by Lemma C.2 of Fan et al. (2011), $|v_{i_1,s\Delta}v_{i_2,s\Delta}|$ still satisfies the exponential-type tail condition,

$$\max_{1 \le i \le d} P\left(|v_{i_1, s\Delta} v_{i_2, s\Delta} - \mathsf{E}[v_{i_1, s\Delta} v_{i_2, s\Delta}]| > c \right) \le \exp(1 - (c/b_4)^{\gamma_5}),\tag{B.2}$$

for $1 \leq i_1, i_2 \leq d$, $s = 1, \dots, n$, c > 0, some b_4 , and $\gamma_5 \in (0, \gamma_4/2)$. Therefore, using the arguments in the proof of Lemma A.3 in Fan et al. (2011), we can show that

$$P\left(\left\|n^{-1}\sum_{s=1}^{n}\boldsymbol{v}_{s\Delta}\boldsymbol{v}_{s\Delta}^{\mathsf{T}}-\boldsymbol{\Sigma}_{v}\right\|_{\mathsf{MAX}}>C_{2}\sqrt{\frac{\log d}{n}}\right)=O\left(\frac{1}{d^{2}}\right)$$

for some positive constant C_2 , which proves part (i). Parts (ii) and (iii) are similar to part (i) and can be obtained from the inequalities derived in Lemma B.1 of Fan et al. (2011). As for part (iv), we have

$$P\left(\left\|n^{-1}\sum_{s=1}^{n}\boldsymbol{\Lambda}_{H}^{^{\mathsf{T}}}\boldsymbol{v}_{s\Delta}\boldsymbol{v}_{s\Delta}^{^{\mathsf{T}}}\boldsymbol{\Lambda}_{H}-\boldsymbol{\Lambda}_{H}^{^{\mathsf{T}}}\boldsymbol{\Sigma}_{v}\boldsymbol{\Lambda}_{H}\right\|_{\mathsf{MAX}} > dm_{v,d}\cdot x\right)$$

$$\leq K_{H}^{2}\max_{1\leq j_{1},j_{2}\leq K_{H}}P\left(\left|n^{-1}\sum_{s=1}^{n}\boldsymbol{\lambda}_{H,j_{1}}^{^{\mathsf{T}}}\boldsymbol{v}_{s\Delta}\boldsymbol{v}_{s\Delta}^{^{\mathsf{T}}}\boldsymbol{\lambda}_{H,j_{2}}-\boldsymbol{\lambda}_{H,j_{1}}^{^{\mathsf{T}}}\boldsymbol{\Sigma}_{v}\boldsymbol{\lambda}_{H,j_{2}}\right| > dm_{v,d}\cdot x\right).$$
(B.3)

Applying similar arguments in (B.1) and (B.2) and using Lemma C.2 of Fan et al. (2011) under Assumption 5(iv), we have

$$\max_{1 \le j_1, j_2 \le K_H} P\left((dm_{v,d})^{-1} | \boldsymbol{\lambda}_{H,j_1}^{^{\mathsf{T}}} \boldsymbol{v}_{s\Delta} \boldsymbol{v}_{s\Delta}^{^{\mathsf{T}}} \boldsymbol{\lambda}_{H,j_2} - \mathsf{E}[\boldsymbol{\lambda}_{H,j_1}^{^{\mathsf{T}}} \boldsymbol{v}_{s\Delta} \boldsymbol{v}_{s\Delta}^{^{\mathsf{T}}} \boldsymbol{\lambda}_{H,j_2}] | > c \right) \le \exp(1 - (c/b_5)^{\gamma_6}), \quad (B.4)$$

for $\gamma_6 \in (0, \gamma_2 \gamma_3/(\gamma_2 + \gamma_3))$, c > 0, and some $b_5 > 0$ which does not depend on n and d. Since $|\boldsymbol{\lambda}_{H,j_1}^{\mathsf{T}} \boldsymbol{v}_{s\Delta} \boldsymbol{v}_{s\Delta}^{\mathsf{T}} \boldsymbol{\lambda}_{H,j_2}|$ satisfies the strong mixing condition, we can follow the same arguments as the proof of Lemma B.1 in Fan et al. (2011) by applying the Bernstein's inequality in Theorem 1 of Merlevède et al. (2011) to obtain

$$P\left((dm_{v,d})^{-1} \left| n^{-1} \sum_{s=1}^{n} \boldsymbol{\lambda}_{H,j_1}^{\mathsf{T}} \boldsymbol{v}_{s\Delta} \boldsymbol{v}_{s\Delta}^{\mathsf{T}} \boldsymbol{\lambda}_{H,j_2} - \boldsymbol{\lambda}_{H,j_1}^{\mathsf{T}} \boldsymbol{\Sigma}_{v} \boldsymbol{\lambda}_{H,j_2} \right| > C_3 \sqrt{\frac{\log d}{n}} \right) = O\left(\frac{1}{d^2}\right), \quad (B.5)$$

for some positive constant C_3 , which only depends on γ_1 , γ_6 and b_5 . Then by (B.3) and (B.5), we can complete the proof of part (iv).

Lemma B.3. Under Assumptions 1, 3 and 6, we have (i) $\|n^{-1/2} \sum_{s=1}^{n} u_{s\Delta} v_{s\Delta}^{\mathsf{T}}\|_{\mathsf{MAX}} = O_P((\log d/n)^{1/2});$ (ii) $\|n^{-1/2} \sum_{s=1}^{n} u_{s\Delta} g_{s\Delta}^{\mathsf{T}}\|_{\mathsf{MAX}} = O_P((\log d/n)^{1/2});$ (iii) $\|n^{-1/2} \sum_{s=1}^{n} v_{s\Delta} f_{s\Delta}^{\mathsf{T}}\|_{\mathsf{MAX}} = O_P((\log d/n)^{1/2});$ (iv) $\|n^{-1/2} \sum_{s=1}^{n} g_{s\Delta} f_{s\Delta}^{\mathsf{T}}\|_{\mathsf{MAX}} = O_P((1/n)^{1/2});$ (v) $\|n^{-1/2} \sum_{s=1}^{n} G_{s\Delta}(n^{-1/2} F_{s\Delta}^{\mathsf{T}})\|_{\mathsf{MAX}} = O_P((1/n)^{1/2}).$

Proof. (i) The proof is similar to that of Lemma 11 in Tao et al. (2013b). Since $\sum_{s=1}^{n} \boldsymbol{u}_{s\Delta} \boldsymbol{v}_{s\Delta}^{\mathsf{T}} = \sum_{s=1}^{n} \boldsymbol{u}_{s\Delta} \boldsymbol{V}_{s\Delta}^{\mathsf{T}} - \sum_{s=1}^{n} \boldsymbol{u}_{s\Delta} \boldsymbol{V}_{(s-1)\Delta}^{\mathsf{T}}$, we only need to prove

$$n^{-1/2} \sum_{s=1}^{n} \boldsymbol{u}_{s\Delta} \boldsymbol{V}_{s\Delta}^{\mathsf{T}} = O_P((\log d/n)^{1/2})$$
(B.6)

and

$$n^{-1/2} \sum_{s=1}^{n} \boldsymbol{u}_{s\Delta} \boldsymbol{V}_{(s-1)\Delta}^{\mathsf{T}} = O_P((\log d/n)^{1/2}).$$
(B.7)

The proofs of (B.6) and (B.7) are similar, so we only provide the former. Denote

$$\Omega_0 = \{ \max_{1 \le i \le d} \max_{1 \le s \le n} |u_{i,s\Delta}| \le 1 \}.$$

Using the Bonferroni and Markov inequalities, we have

$$\mathsf{P}(\Omega_0^c) = \mathsf{P}\left(\max_{1 \le i \le d} \max_{1 \le s \le n} |u_{i,s\Delta}| > 1\right) \le n de^{C_\sigma/2 - n}.$$
(B.8)

Note that U_t and V_t are independent. Conditional on the whole path of U_t , we have

$$\begin{aligned}
\mathsf{P}\left(\left\|n^{-1/2}\sum_{s=1}^{n}\boldsymbol{u}_{s\Delta}\boldsymbol{V}_{s\Delta}^{\mathsf{T}}\right\|_{\mathsf{MAX}} > c(\log d/n)^{1/2}\right) \\
\leq \mathsf{P}\left(\left\|n^{-1/2}\sum_{s=1}^{n}\boldsymbol{u}_{s\Delta}\boldsymbol{V}_{s\Delta}^{\mathsf{T}}\right\|_{\mathsf{MAX}} > c(\log d/n)^{1/2},\Omega_{0}\right) + \mathsf{P}\left(\Omega_{0}^{c}\right) \\
\leq \mathsf{E}\left[\mathsf{P}\left(\left\|n^{-1/2}\sum_{s=1}^{n}\boldsymbol{u}_{s\Delta}\boldsymbol{V}_{s\Delta}^{\mathsf{T}}\right\|_{\mathsf{MAX}} > c(\log d/n)^{1/2},\Omega_{0}\right|\boldsymbol{U}_{t},\,t\in[0,1]\right)\right] + O(nde^{-n}) \\
= \mathsf{E}\left[\mathsf{P}\left(\left\|n^{-1/2}\sum_{s=1}^{n}\boldsymbol{u}_{s\Delta}\boldsymbol{V}_{s\Delta}^{\mathsf{T}}\right\|_{\mathsf{MAX}} > c(\log d/n)^{1/2}\left|\Omega_{0},\boldsymbol{U}_{t},\,t\in[0,1]\right)\right] + O(nde^{-n}). \quad (B.9)
\end{aligned}$$

Note that conditional on the path of U_t and Ω_0 , $u_{s\Delta}V_{s\Delta}^{\dagger}$ satisfies the same mixing condition and exponential-tail condition for $V_{s\Delta}$, and the coefficients in these conditions only depend on γ_1, γ_2 , and b_1 . Thus we can apply the Bernstein's inequality in Theorem 1 of Merlevède et al. (2011) to obtain (letting $\bar{c} = c(\log d/n)^{1/2}$)

$$\mathsf{P}\left(\left\|n^{-1/2}\sum_{s=1}^{n}\boldsymbol{u}_{s\Delta}\boldsymbol{V}_{s\Delta}^{\mathsf{T}}\right\|_{\mathsf{MAX}} > \bar{c}\left|\Omega_{0},\boldsymbol{U}_{t},\,t\in[0,1]\right) \\ \leq nd^{2}\exp\left(-\frac{\bar{c}^{\gamma}}{C_{4}}\right) + d^{2}\exp\left(-\frac{\bar{c}^{2}}{C_{5}(1+C_{6}n)}\right) + d^{2}\exp\left(-\frac{\bar{c}^{2}}{C_{7}n}\exp\left(\frac{\bar{c}^{\gamma(1-\gamma)}}{C_{8}((\log\bar{c})^{\gamma})}\right)\right) \\ = O(1/d^{2}), \tag{B.10}$$

when $(\log d)^{2/\gamma-1} = o(n)$ and c is large enough, where $\gamma = 1/\gamma_1 + 1/\gamma_2$, and C_4-C_8 only depends on γ_1, γ_2 , and b_1 . Therefore (B.10) holds true uniformly for all path of U_t satisfying Ω_0 . Combining (B.9) and (B.10), we can prove (B.6). The proofs of (ii)–(v) are similar to that of (i) by choosing proper \bar{c} . So we omit them to save space.

Lemma B.4. Under Assumptions 1–6, we have
(i)
$$\| \boldsymbol{w}^{\mathsf{T}} \boldsymbol{w} - \boldsymbol{\Sigma}_{w} \|_{\mathsf{MAX}} = O_{P}((\log d/n)^{1/2} \cdot n^{2\bar{\tau}_{V}^{*+}});$$

(ii) $\| \boldsymbol{w}^{\mathsf{T}} \boldsymbol{h} \mathbf{D}_{H} \|_{\mathsf{MAX}} = O_{P}\left((\log d/n)^{1/2} \cdot n^{\bar{\tau}_{V}^{*+} + \bar{\tau}_{G}^{*+}} \right);$
(iii) $\| \mathbf{D}_{H} \boldsymbol{h}^{\mathsf{T}} \boldsymbol{w} \|_{1} = O_{P}\left((\log d/n)^{1/2} \cdot n^{\bar{\tau}_{V}^{*+} + \bar{\tau}_{G}^{*+}} \right);$
(iv) $\| \boldsymbol{w} \|_{\mathsf{sp}} = O_{P}(d^{1/2}(\log d/n)^{1/4} \cdot n^{\bar{\tau}_{V}^{*+}} + m^{1/2}_{U,d} + n^{\bar{\tau}_{V}} m^{1/2}_{v,d});$
(v) $\| \mathbf{\Lambda}_{H}^{\mathsf{T}} \boldsymbol{w}^{\mathsf{T}} \|_{\mathsf{sp}} = O_{P}(d^{1/2}(\log d/n)^{1/4} \cdot n^{\bar{\tau}_{V}^{*+}} + d^{1/2}(m^{1/2}_{U,d} + n^{\bar{\tau}_{V}} m^{1/2}_{v,d}));$
where $\bar{\tau}_{G}^{*+} = (1/2 + \bar{\tau}_{G})_{+}, \ \bar{\tau}_{V}^{*} = (1/2 + \bar{\tau}_{V}), \ and \ \bar{\tau}_{V}^{*+} = (1/2 + \bar{\tau}_{V})_{+}.$

Proof. For part (i), recall that $ww^{\mathsf{T}} = \sum_{s=1}^{n} (\boldsymbol{u}_{s\Delta} + \mathbf{D}_{V}\boldsymbol{v}_{s\Delta})^{\mathsf{T}} (\boldsymbol{u}_{s\Delta} + \mathbf{D}_{V}\boldsymbol{v}_{s\Delta})$ and $\boldsymbol{\Sigma}_{w} = \boldsymbol{\Sigma}_{U} + n\mathbf{D}_{V}\boldsymbol{\Sigma}_{v}\mathbf{D}_{V}$. By Lemmas B.1(i), B.2(i) and B.3(i), we have

$$\begin{aligned} \left\| \boldsymbol{w}^{\mathsf{T}} \boldsymbol{w} - \boldsymbol{\Sigma}_{\boldsymbol{w}} \right\|_{\mathsf{MAX}} &\leq \left\| \sum_{s=1}^{n} \boldsymbol{u}_{s\Delta} \boldsymbol{u}_{s\Delta}^{\mathsf{T}} - \boldsymbol{\Sigma}_{\boldsymbol{U}} \right\|_{\mathsf{MAX}} + 2n^{1/2} \|\mathbf{D}_{\boldsymbol{V}}\| \left\| n^{-1/2} \sum_{s=1}^{n} \boldsymbol{u}_{s\Delta} \boldsymbol{v}_{s\Delta}^{\mathsf{T}} \right\|_{\mathsf{MAX}} \\ &+ n \|\mathbf{D}_{\boldsymbol{V}}\|_{\mathsf{sp}}^{2} \left\| n^{-1} \sum_{s=1}^{n} \boldsymbol{v}_{s\Delta} \boldsymbol{v}_{s\Delta}^{\mathsf{T}} - \boldsymbol{\Sigma}_{\boldsymbol{v}} \right\|_{\mathsf{MAX}} \\ &= O_{P} \left((\log d/n)^{1/2} \cdot n^{2\overline{\tau}_{\boldsymbol{V}}^{*+}} \right). \end{aligned}$$

For part (ii), by Lemmas B.1(iii), B.2(iii), B.3(ii) and B.3(iii), we have

$$\begin{split} \left\| \boldsymbol{w}^{\mathsf{T}} \boldsymbol{\hbar} \mathbf{D}_{H} \right\|_{\mathsf{MAX}} &\leq \\ \left\| \sum_{s=1}^{n} \boldsymbol{u}_{s\Delta} \boldsymbol{f}_{s\Delta}^{\mathsf{T}} \right\|_{\mathsf{MAX}} + n \| \mathbf{D}_{V} \|_{\mathsf{sp}} \left\| n^{-1} \sum_{s=1}^{n} \boldsymbol{v}_{s\Delta} \boldsymbol{g}_{s\Delta}^{\mathsf{T}} \right\|_{\mathsf{MAX}} \| \mathbf{D}_{G} \|_{\mathsf{sp}} \\ &+ n^{1/2} \left\| n^{-1/2} \sum_{s=1}^{n} \boldsymbol{u}_{s\Delta} \boldsymbol{g}_{s\Delta}^{\mathsf{T}} \right\|_{\mathsf{MAX}} \| \mathbf{D}_{G} \|_{\mathsf{sp}} + \| \mathbf{D}_{V} \|_{\mathsf{sp}} \left\| \sum_{s=1}^{n} \boldsymbol{v}_{s\Delta} \boldsymbol{f}_{s\Delta}^{\mathsf{T}} \right\|_{\mathsf{MAX}} \\ &= O_{P} \left((\log d/n)^{1/2} \cdot (1 + n^{1+\bar{\tau}_{V}+\bar{\tau}_{G}} + n^{1/2+\bar{\tau}_{G}} + n^{1/2+\bar{\tau}_{V}}) \right) \\ &= O_{P} \left((\log d/n)^{1/2} \cdot (1 + n^{1/2+\bar{\tau}_{V}})(1 + n^{1/2+\bar{\tau}_{G}}) \right) \\ &= O_{P} \left((\log d/n)^{1/2} \cdot n^{\bar{\tau}_{V}^{*+} + \bar{\tau}_{G}^{*+}} \right). \end{split}$$

Part (iii) follows from part (ii) as $\| \varkappa^{\mathsf{T}} w \|_{1} \leq K_{H} \| \varkappa^{\mathsf{T}} w \|_{\mathsf{MAX}}$. Part (iv) follows from $\| w \|_{\mathsf{sp}} = \| w^{\mathsf{T}} w \|_{\mathsf{sp}}^{1/2} \leq (d \| w^{\mathsf{T}} w - \Sigma_{w} \|_{\mathsf{MAX}} + \| \Sigma_{w} \|_{\mathsf{sp}})^{1/2}$ and $\| \Sigma_{w} \|_{\mathsf{sp}} \leq (\| \Sigma_{w} \|_{1} \| \Sigma_{w} \|_{\infty})^{1/2} = \| \Sigma_{w} \|_{1} = O(m_{U,d} + n^{2\tau_{V}^{*}} m_{v,d})$. Lastly, we consider part (v). By Lemmas B.1(iv) and B.2(iv), we have

$$\begin{split} \|\mathbf{\Lambda}_{H}^{'}\boldsymbol{w}^{'}\|_{\mathsf{sp}} &= \|\mathbf{\Lambda}_{H}^{'}\boldsymbol{w}^{'}\boldsymbol{w}\mathbf{\Lambda}_{H}\|_{\mathsf{sp}}^{1/2} \\ &\leq (\|\mathbf{\Lambda}_{H}^{\mathsf{T}}(\boldsymbol{w}^{\mathsf{T}}\boldsymbol{w}-\boldsymbol{\Sigma}_{w})\mathbf{\Lambda}_{H}\|_{\mathsf{sp}} + \|\mathbf{\Lambda}_{H}^{\mathsf{T}}\boldsymbol{\Sigma}_{w}\mathbf{\Lambda}_{H}\|_{1})^{1/2} \\ &\leq (K_{H}\|\mathbf{\Lambda}_{H}^{\mathsf{T}}(\boldsymbol{w}^{\mathsf{T}}\boldsymbol{w}-\boldsymbol{\Sigma}_{w})\mathbf{\Lambda}_{H}\|_{\mathsf{MAX}} + \|\mathbf{\Lambda}_{H}^{\mathsf{T}}\|_{1}\|\mathbf{\Sigma}_{w}\|_{1}\|\mathbf{\Lambda}_{H}\|_{1})^{1/2} \\ &= O_{P}\left(d^{1/2}(m_{U,d}^{1/2}+n^{\bar{\tau}_{V}^{*}}m_{v,d}^{1/2})(\log d/n)^{1/4} \cdot n^{\bar{\tau}_{V}^{*+}} + d^{1/2}(m_{U,d}^{1/2}+n^{\bar{\tau}_{V}^{*}}m_{v,d}^{1/2})\right) \\ &= O_{P}\left(d^{1/2}(m_{U,d}^{1/2}+n^{\bar{\tau}_{V}^{*}}m_{v,d}^{1/2})(1+n^{\bar{\tau}_{V}^{*+}}(\log d/n)^{1/4})\right). \end{split}$$

Lemma B.5. Under Assumptions 1, 3 and 6, we have (i) $\|n^{-1}\sum_{s=1}^{n} (n^{-1/2}\boldsymbol{G}_{s\Delta}^{c})(n^{-1/2}\boldsymbol{G}_{s\Delta}^{c^{\mathsf{T}}})\|_{\mathsf{MAX}} = O_P(n^{-1});$ (ii) $n^{-1}\sum_{s=1}^{n} \boldsymbol{F}_{s\Delta}^{c} \boldsymbol{F}_{s\Delta}^{c^{\mathsf{T}}} \xrightarrow{d} \int_{0}^{1} (\int_{0}^{t} \boldsymbol{\sigma}_{fu} \mathrm{d}\boldsymbol{B}_{u}^{F} - \int_{0}^{1} \boldsymbol{\sigma}_{fu} \mathrm{d}\boldsymbol{B}_{u}^{F}) (\int_{0}^{t} \boldsymbol{\sigma}_{fu} \mathrm{d}\boldsymbol{B}_{u}^{F} - \int_{0}^{1} \boldsymbol{\sigma}_{fu} \mathrm{d}\boldsymbol{B}_{u}^{F})^{\mathsf{T}} \mathrm{d}t;$ (iii) $\|n^{-1}\sum_{s=1}^{n} (n^{-1/2}\boldsymbol{G}_{s\Delta}^{c})\boldsymbol{F}_{s\Delta}^{c^{\mathsf{T}}}\|_{\mathsf{MAX}} = O_P(n^{-1});$ where $\boldsymbol{G}_{s\Delta}^{c} = \boldsymbol{G}_{s\Delta} - \overline{\boldsymbol{G}}, \ \boldsymbol{F}_{s\Delta}^{c} = \boldsymbol{F}_{s\Delta} - \overline{\boldsymbol{F}}, \ \overline{\boldsymbol{G}} = n^{-1}\sum_{s=1}^{n} \boldsymbol{G}_{s\Delta}, \ and \ \overline{\boldsymbol{F}} = n^{-1}\sum_{s=1}^{n} \boldsymbol{F}_{s\Delta}.$

Proof. Parts (i) and (ii) are trivial. For part (iii), by Lemma B.3(v), we have

$$\left\| n^{-3/2} \sum_{s=1}^{n} \boldsymbol{G}_{s\Delta}^{c} \boldsymbol{F}_{s\Delta}^{c^{\mathsf{T}}} \right\|_{\mathsf{MAX}} = n^{-1/2} \left\| n^{-1} \sum_{s=1}^{n} \boldsymbol{G}_{s\Delta} \boldsymbol{F}_{s\Delta}^{\mathsf{T}} - \overline{\boldsymbol{G}} \ \overline{\boldsymbol{F}}^{\mathsf{T}} \right\|_{\mathsf{MAX}} = O_P(n^{-1}).$$

Lemma B.6. Under Assumptions 1–8, we have (i) $\|n^{-1}\mathscr{H}^{*c^{\mathsf{T}}}\mathscr{H}^{*c}\|_{\mathsf{sp}} = O_P(1);$ (ii) $\|\widehat{\mathscr{H}}^{*c}\Xi^{\mathsf{T}} - \mathscr{H}^c\|_{\mathsf{MAX}} \leq \|\widehat{\mathscr{H}}^{*c}\Xi^{\mathsf{T}} - \mathscr{H}^c\|_{\mathsf{sp}} = O_P\left(n^{-2\underline{\tau}_G^{*-}} \cdot a_{nd}\right);$ (iii) $n^{-1}\|\Xi\widehat{\mathscr{H}}^{*c^{\mathsf{T}}}\widehat{\mathscr{H}}^{*c}\Xi^{\mathsf{T}} - \mathscr{H}^{*c^{\mathsf{T}}}\mathscr{H}^{*c}\| = O_P\left(n^{-1/2-2\underline{\tau}_G^{*-}} \cdot a_{nd}\right);$ where Ξ is defined in Lemma 3.1.

Proof. (i) By Lemma B.4, the dominate term is $n^{-1} \sum_{s=1}^{n} \mathbf{F}_{s\Delta}^{c} \mathbf{F}_{s\Delta}^{c\dagger}$, which is of order $O_{P}(1)$.(ii) By Theorem 3.1, we have

$$\left\|\widehat{\mathscr{H}^*}\Xi^{\mathsf{T}} - (\mathscr{H} - \mathbf{1}_n \boldsymbol{h}_0^{\mathsf{T}})\right\|_{\mathsf{MAX}} \leq \left\|\widehat{\mathscr{H}^*}\Xi^{\mathsf{T}} - (\mathscr{H} - \mathbf{1}_n \boldsymbol{h}_0^{\mathsf{T}})\right\|_{\mathsf{sp}} = O_P\left(n^{-2\underline{\tau}_G^{*-}} \cdot a_{nd}\right).$$
$$\Xi^{\mathsf{T}} - (\mathbf{I}_{\mathsf{MAX}} - \mathbf{1}_{\mathsf{MAX}} \cdot \mathbf{1}_{\mathsf{MAX}}) - (\mathbf{I}_{\mathsf{MAX}} - \mathbf{1}_{\mathsf{MAX}} \cdot \mathbf{1}_{\mathsf{MAX}} \cdot \mathbf{1}_{\mathsf{MAX}}) - \mathbf{1}_{\mathsf{MAX}} \cdot \mathbf{1}_{\mathsf{$$

Since $\widehat{\mathscr{H}}^{*c} \Xi^{\mathsf{T}} = (\mathbf{I}_n - \mathbf{1}_n \mathbf{1}_n^{\mathsf{T}}/n) \widehat{\mathscr{H}} \Xi^{\mathsf{T}}, \ \mathscr{H}^c \Xi^{\mathsf{T}} = (\mathbf{I}_n - \mathbf{1}_n \mathbf{1}_n^{\mathsf{T}}/n) (\mathscr{H} - \mathbf{1}_n \mathbf{h}_0^{\mathsf{T}}) \text{ and } \|\mathbf{I}_n - \mathbf{1}_n \mathbf{1}_n^{\mathsf{T}}/n\|_{\mathsf{sp}} = 1, \text{ we then have } \|\widehat{\mathscr{H}}^{*c} \Xi^{\mathsf{T}} - \mathscr{H}^c\|_{\mathsf{sp}} = O_P \left(n^{-2\underline{\tau}_G^{*-}} \cdot a_{nd}\right).$

(iii) The result follows by noticing that

$$\begin{aligned} & \left\| \Xi \widehat{\mathscr{H}}^{*c^{\mathsf{T}}} \widehat{\mathscr{H}}^{*c} \Xi^{\mathsf{T}} - \mathscr{H}^{*c^{\mathsf{T}}} \mathscr{H}^{*c} \right\| \\ & \leq \left\| \widehat{\mathscr{H}}^{*c} \Xi^{\mathsf{T}} - \mathscr{H}^{c} \right\|_{\mathsf{sp}} \left\| \mathscr{H}^{c} \right\|_{\mathsf{sp}} + \left\| \widehat{\mathscr{H}}^{*c} \Xi^{\mathsf{T}} - \mathscr{H}^{c} \right\|_{\mathsf{sp}}^{2} \\ & = O_{P} \left(n^{1/2 - 2\underline{\tau}_{G}^{*-}} \cdot a_{nd} \right). \end{aligned}$$

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