# Latent Factor Analysis in Short Panels 

Alain-Philippe Fortin ${ }^{1,2}$, Patrick Gagliardini ${ }^{3,2}$, Olivier Scaillet ${ }^{1,2 *}$

September 13, 2023


#### Abstract

We develop inferential tools for latent factor analysis in short panels. The pseudo maximum likelihood setting under a large cross-sectional dimension $n$ and a fixed time series dimension $T$ relies on a diagonal $T \times T$ covariance matrix of the errors without imposing sphericity or Gaussianity. We outline the asymptotic distributions of the latent factor and error covariance estimates as well as of an asymptotically uniformly most powerful invariant (AUMPI) test based on the likelihood ratio statistic for tests of the number of factors. We derive the AUMPI characterization from inequalities ensuring the monotone likelihood ratio property for positive definite quadratic forms in normal variables. An empirical application to a large panel of monthly U.S. stock returns separates date after date systematic and idiosyncratic risks in short subperiods of bear vs. bull market based on the selected number of factors. We observe an uptrend in idiosyncratic volatility while the systematic risk explains a large part of the cross-sectional total variance in bear markets but is not driven by a single factor. Rank tests reveal that observed factors struggle spanning latent factors with a discrepancy between the dimensions of the two factor spaces decreasing over time.


Keywords: Latent factor analysis, uniformly most powerful invariant test, panel data, large $n$ and fixed $T$ asymptotics, equity returns. JEL codes: C12, C23, C38, C58, G12.

[^0]
## 1 Introduction

Latent variable models have been used for a long time in econometrics (Aigner et al. (1984)). Here, we study large cross-sectional latent factor models with small time dimension. Two common methods for estimation of latent factor spaces are principal component analysis (PCA) and factor analysis (FA), see Anderson (2003) Chapters 11 and 14. They cover multiple applications in finance and economics as well as in social sciences in general. They are often used in exploratory analysis of data. In recent work, Fortin, Gagliardini and Scaillet (FGS, 2022) show how we can use PCA to conduct inference on the number of factors in such models without making Gaussian assumptions. Their methodology relies on sphericity of the idiosyncratic variances since this restriction is both necessary and sufficient for consistency of latent factor estimates with small $T$ (Theorem 4 of Bai (2003)). In PCA, sphericity allows to identify the number $k$ of factors from the $k$ first eigenvalue spacings being larger than zero, and being zero the subsequent ones. On the contrary, the FA strategy does not exploit eigenvalue spacings and does not require sphericity. However, inference with small $T$ up to now mostly relies on (often restrictive) assumptions such as Gaussian variables (with a notable exception by Anderson and Amemiya (1988)) and error homoschedasticity across sample units. Those are untenable assumptions in our application with stock returns.

A central and practical issue in applied work with latent factors is to determine the number of factors. For models with unobservable (latent) factors only, Connor and Korajczyk (1993) are the first to develop a test for the number of factors for large balanced panels of individual stock returns in time-invariant models under covariance stationarity and homoskedasticity. Unobservable factors are estimated by the method of asymptotic principal components developed by Connor and Korajczyk (1986) (see also Stock and Watson (2002)). For heteroskedastic settings, the recent literature on large balanced panels with static factors has extended the toolkit available to researchers. A first strand of that literature focuses on consistent estimation procedures for the number of factors. Bai and Ng (2002) introduce a penalized least-squares strategy to estimate the number of factors,
at least one. Ando and Bai (2015) extend that approach when explanatory variables are present in the linear specification (see Bai (2009) for homogeneous regression coefficients). Onatski (2010) looks at the behavior of differences in adjacent eigenvalues to determine the number of factors when $n$ and $T$ are both large and comparable. Ahn and Horenstein (2013) opt for a similar strategy based on eigenvalue ratios. Caner and Han (2014) propose an estimator with a group bridge penalization to determine the number of unobservable factors. Based on the framework of Gagliardini, Ossola and Scaillet (2016), Gagliardini, Ossola and Scaillet (2019) build a simple diagnostic criterion for approximate factor structure in large panel datasets. Given observable factors, the criterion checks whether the errors are weakly cross-sectionally correlated, or share one or more unobservable common factors (interactive effects), and selects their number; see Gagliardini, Ossola and Scaillet (2020) for a survey of estimation of large dimensional conditional factor models in finance. A second strand of that literature develops inference procedures for hypotheses on the number of latent factors. Onatski (2009) deploys a characterization of the largest eigenvalues of a Wishart-distributed covariance matrix with large dimensions in terms of the Tracy-Widom Law. To get a Wishart distribution, Onatski (2009) assumes either Gaussian errors, or $T$ much larger than $n$. Kapetanios (2010) uses subsampling to estimate the limit distribution of the adjacent eigenvalues.

This paper puts forward methodological and empirical contributions that complement the above literature. (i) On the methodological side, we extend the inferential tools of FA to non-Gaussian and non-i.i.d. settings. First, we characterize the asymptotic distribution of FA estimators obtained under a pseudo maximum likelihood approach where the time-series dimension is held fixed while the cross-sectional dimension diverges. Hence, the asymptotic analysis targets short panels, and allows for cross-sectionally heteroschedastic and weakly dependent errors. Cochrane (2005, p. 226) argues in favour of the development of appropriate large- $n$ small- $T$ tools for evaluating asset pricing models, a problem only partially addressed in finance. In a short panel setting, Zaffaroni (2019) considers inference for latent factors in conditional linear asset pricing models under sphericity based on PCA, including estimation of the number of factors. ${ }^{1}$ The small $T$ setting

[^1]mitigates concerns for panel unbalancedness and corresponds to a locally time-invariant factor structure accommodating globally time-dependent features of general forms. It is also appealing to macroeconomic data observed quarterly. For the sake of space, we put part of the theory, namely inference for FA estimates, in the Online Appendix (OA). We refer to Bai and Li (2016) for inference when $n$ and $T$ are both large (see Bai and Li (2012) for the cross-sectional independent case). Second, we use our new theoretical results for FA to develop testing procedures for the number of latent factors in a short panel which rely neither on sphericity nor Gaussianity, thereby extending tests based on eigenvalues, as in Onatski (2009), to small $T$, and as in FGS, to nonspherical errors, thanks to an FA device. We derive the Asymptotically Uniformly Most Powerful Invariant (AUMPI) property of the FA likelihood ratio (LR) test statistic in the non-Gaussian case under inequality restrictions on the DGP parameters, and cover inference with weak factors. The AUMPI property is rare and sought-after in testing procedures (see Engle (1984) for a discussion), and often holds only under restrictive assumptions such as Gaussianity. (ii) On the empirical side, we apply our FA methodology to panels of monthly U.S. stock returns with large cross-sectional and small time-series dimensions, and investigate how the number of driving factors changes over time and particular periods. Furthermore, date after date, we provide a novel separation of the risk coming from the systematic part and the risk coming from the idiosyncratic part of returns in short subperiods of bear vs. bull market based on the selected number of factors. We observe an uptrend in idiosyncratic volatility (see also Campbell et al. (2023)) while the systematic risk explains a large part of the cross-sectional total variance in bear markets but is not driven by a single factor. We also investigate whether standard observed factors span the estimated latent factors using rank tests in order to suit our fixed $T$ setting. Observed factors struggle spanning latent factors with a discrepancy between the dimensions of the two factor spaces decreasing over time.
the ex-post risk premia (Shanken (1992)) associated to observable factors (see Kleibergen and Zhan (2023) for robustidentification inference based a continuous updating generalized method of moments). Kim and Skoulakis (2018) deals with the error-in-variable problem of the two-pass methodology with small $T$ by regression-calibration under sphericity and a block-dependence structure.

The outline of the paper is as follows. In Section 2, we consider a linear latent factor model and introduce test statistics on the number of latent factors based on FA. Section 3 presents the asymptotic distributional theory for inference in short panels under a block-dependence structure to allow for weak dependence in the cross-section. Section 4 discusses three special cases, i.e. Gaussian errors, settings where the asymptotic distribution under Gaussian errors still holds for the test statistics, and spherical errors. Section 5 is dedicated to local asymptotic power and AUMPI tests. We provide our empirical application in Section 6 and our concluding remarks in Section 7. Appendices A and B gather the regularity assumptions and proofs of the main theoretical results. We place all omitted proofs and additional analyses in Appendices C-E in Online Appendix (OA). Besides, we gather all explicit formulas not listed in the core text but useful for coding in an online "Supplementary Materials for Coding" (SMC) attached to the replication files. We also put there other numerical checks and a Monte Carlo assessment of size and power for the LR test statistic.

## 2 Test statistics based on Factor Analysis

We consider the linear Factor Analysis (FA) model (e.g. Anderson (2003)):

$$
\begin{equation*}
y_{i}=\mu+F \beta_{i}+\varepsilon_{i}, \quad i=1, \ldots, n \tag{1}
\end{equation*}
$$

where $y_{i}=\left(y_{i, 1}, \ldots, y_{i, T}\right)^{\prime}$ and $\varepsilon_{i}=\left(\varepsilon_{i, 1}, \ldots, \varepsilon_{i, T}\right)^{\prime}$ are $T$-dimensional vectors of observed data and unobserved error terms for individual $i$. The $k$-dimensional vectors $\beta_{i}=\left(\beta_{i, 1}, \ldots, \beta_{i, k}\right)^{\prime}$ are latent individual effects, while $\mu$ and $F$ are a $T \times 1$ vector and a $T \times k$ matrix of unknown parameters. The number of latent factors $k$ is an unknown integer smaller than $T$. In matrix notation, model (1) reads $Y=\mu 1_{n}^{\prime}+F \beta^{\prime}+\varepsilon$, where $Y$ and $\varepsilon$ are $T \times n$ matrices, $\beta$ is the $n \times k$ matrix with rows $\beta_{i}^{\prime}$, and $1_{n}$ is a $n$-dimensional vector of ones.

Assumption 1 The $T \times T$ matrix $V_{\varepsilon}=\lim _{n \rightarrow \infty} E\left[\frac{1}{n} \varepsilon \varepsilon^{\prime}\right]$ is diagonal.
Matrix $V_{\varepsilon}$ is the limit cross-sectional average of the - possibly heterogeneous - errors' unconditional
variance-covariance matrices. The diagonality condition in Assumption 1 is standard in FA (in the more restrictive formulation involving i.i.d. data).

In our empirics with a large cross-sectional panel of returns for $n$ assets over a short time span with $T$ periods, vectors $y_{i}$ and $\varepsilon_{i}$ stack the monthly returns and the idiosyncratic errors of stock $i$. Any row vector $f_{t}^{\prime}:=\left(f_{t, 1}, \ldots, f_{t, k}\right)$ of matrix $F$ yields the latent factor values in a given month $t$, and vector $\beta_{i}$ collects the factor loadings of stock $i$. In our finance application, we assume the No-Arbitrage (NA) principle to hold, so that the entries $\mu_{t}$ in the intercept vector in Equation (1) account for the (possibly time-varying) risk-free rate and (possibly non-zero) cross-sectional mean of stock betas. ${ }^{2}$ Thus, the linear FA model (1) yields $y_{i, t}=\mu_{t}+f_{t}^{\prime} \beta_{i}+\varepsilon_{i, t}$, that is the standard formulation in asset pricing. We cover the Capital Asset Pricing Model (CAPM) when the single latent factor is the excess return of the market portfolio. Assumption 1 allows for serial dependence in idiosyncratic errors in the form of martingale difference sequences, like individual GARCH and Stochastic Volatility (SV) processes, as well as weak cross-sectional dependence (see Assumption 2 below). It also accommodates common time-varying components in idiosyncratic volatilities by allowing different entries along the diagonal of $V_{\varepsilon}$; see Renault, Van Der Heijden and Werker (2022) for arbitrage pricing in such settings. ${ }^{3}$

This paper focuses mainly on testing hypotheses on the number of latent factors $k$ when $T$ is

[^2]fixed and $n \rightarrow \infty$. The fixed $T$ perspective makes FA especially well-suited for applications with short panels. Indeed, we work conditionally on the realizations of the latent factors $F$ and treat their values as parameters to estimate. In comparison with the standard small $n$ and large $T$ framework in traditional asset pricing (see e.g. Shanken (1992) with observable factors), here factors and loadings are interchanged in the sense that the $\beta_{i}$ and $F$ play the roles of the "factors" and the "factor loadings" in FA. We depart from classical FA since the $\beta_{i}$ are not considered as random effects (e.g. with a Gaussian distribution) but rather as fixed effects, namely incidental parameters. ${ }^{4}$ Moreover, in Assumption 1, we neither assume Gaussianity nor we impose sphericity of the covariance matrix of the error terms. Besides we accommodate weak cross-sectional dependence and ARCH effects in idiosyncratic errors (see Section 3). Hence, the FA estimators defined below correspond to maximizers of a Gaussian pseudo likelihood. By-products of our analysis are the feasible asymptotic distributions of FA estimators of $F$ and $V_{\varepsilon}$ in more general settings than in the available literature (e.g. Anderson and Amemiya (1988)), which we present in Appendix D.

The test statistics we consider for conducting inference on the number of latent factors $k$ are functions of the elements of the symmetric matrix

$$
\begin{equation*}
\hat{S}=\hat{V}_{\varepsilon}^{-1 / 2} M_{\hat{F}, \hat{V}_{\varepsilon}}\left(\hat{V}_{y}-\hat{V}_{\varepsilon}\right) M_{\hat{F}, \hat{V}_{\varepsilon}}^{\prime} \hat{V}_{\varepsilon}^{-1 / 2} \tag{2}
\end{equation*}
$$

where $\hat{V}_{y}=\frac{1}{n} \tilde{Y} \tilde{Y}^{\prime}$ is the sample (cross-sectional) variance matrix (the $n$ columns of $\tilde{Y}$ are $y_{i}-\bar{y}$ and $\bar{y}=\frac{1}{n} \sum_{i=1}^{n} y_{i}$ is the vector of cross-sectional means), $M_{F, V}:=I_{T}-F\left(F^{\prime} V^{-1} F\right)^{-1} F^{\prime} V^{-1}$ is the Generalized Least Squares (GLS) projection matrix orthogonal to $F$ for variance $V$, and $\hat{F}$ and $\hat{V}_{\varepsilon}$ are the FA estimators computed under the assumption that there are $k$ latent factors. In the following, we use the same notation for the matrix-to-vector diag operator and the vector-to-matrix diag operator. Hence, $\operatorname{diag}(A)$ for a matrix $A$ denotes the vector in which we stack the diagonal elements of matrix $A$, and $\operatorname{diag}(a)$ for a vector $a$ denotes a diagonal matrix with the elements of $a$ on the diagonal. From Anderson (2003) Chapter 14, the FA estimators $\hat{F}, \hat{V}_{\varepsilon}$ maximize a Gaussian

[^3]pseudo likelihood (Appendix D.1) and meet the first order conditions: ${ }^{5}$
(FA1) $\quad \operatorname{diag}\left(\hat{V}_{y}\right)=\operatorname{diag}\left(\hat{F} \hat{F}^{\prime}+\hat{V}_{\varepsilon}\right)$, and
(FA2) $\quad \hat{F}$ is the $T \times k$ matrix of eigenvectors of $\hat{V}_{y} \hat{V}_{\varepsilon}^{-1}$ associated to the $k$ largest eigenvalues $1+\hat{\gamma}_{j}, j=1, \ldots, k$, normalized such that $\hat{F}^{\prime} \hat{V}_{\varepsilon}^{-1} \hat{F}=\operatorname{diag}\left(\hat{\gamma}_{1}, \ldots, \hat{\gamma}_{k}\right)$.
The number of degrees of freedom is $d f=\frac{1}{2}\left((T-k)^{2}-T-k\right)$ and it is required that $d f \geq 0 .{ }^{6}$
Statistic $\hat{S}$ in Equation (2) checks if the difference between the sample variance-covariance $\hat{V}_{y}$ and diagonal matrix $\hat{V}_{\varepsilon}$ is a symmetric matrix of reduced rank $k$, with range spanned by the range of $\hat{F}$. The probability limit of $\hat{S}$ is nil under the null hypothesis of $k$ latent factors. We get further insights from the next result.

Proposition 1 Under Assumption 1, (a) the eigenvalues of matrix $\hat{S}$ are: $\hat{\gamma}_{j}$, for $j=k+1, \ldots, T$, and 0 , with multiplicity $k$, where $1+\hat{\gamma}_{j}$ for $j=k+1, \ldots, T$ are the $T-k$ smallest eigenvalues of $\hat{V}_{y} \hat{V}_{\varepsilon}^{-1}$, (b) the squared Frobenius norm is $\|\hat{S}\|^{2}=\sum_{j=k+1}^{T} \hat{\gamma}_{j}^{2}$, (c) $\operatorname{diag}(\hat{S})=0$, and (d) we get

$$
\begin{equation*}
\hat{S}=\hat{V}_{\varepsilon}^{-1 / 2}\left(\frac{1}{n} \hat{\varepsilon} \hat{\varepsilon}^{\prime}\right) \hat{V}_{\varepsilon}^{-1 / 2}-\hat{V}_{\varepsilon}^{-1 / 2} M_{\hat{F}, \hat{V}_{\varepsilon}} \hat{V}_{\varepsilon}^{1 / 2} \tag{3}
\end{equation*}
$$

where $\hat{\varepsilon}=M_{\hat{F}, \hat{V}_{\varepsilon}} \tilde{Y}$ is the $T \times n$ matrix of GLS residuals.
From Proposition 1 (a)-(c), the squared Frobenius norm of matrix $\hat{S}$ multiplied by $n / 2$ coincides at second order with the classical LR statistic in FA, i.e., $L R(k)=-n \sum_{j=k+1}^{T} \log \left(1+\hat{\gamma}_{j}\right)$. ${ }^{7}$ Moreover, from (d) we can interpret matrix $\hat{S}$ in terms of scaled cross-sectional averages of

[^4]squared and cross-products of GLS residuals. In (3), we subtract $\hat{V}_{\varepsilon}^{-1 / 2} M_{\hat{F}, \hat{V}_{\varepsilon}} \hat{V}_{\varepsilon}^{1 / 2}$ and not the identity because the residuals are orthogonal to $\hat{F}$ by construction. From Proposition 1 (c), the diagonal elements of matrix $\hat{S}$ vanish. Those elements are not informative for inference on the number of factors, and can be ignored when constructing the test statistics. This finding is natural because we expect that only the out-of-diagonal elements of $\frac{1}{n} \hat{\varepsilon} \hat{\varepsilon}^{\prime}$, i.e., the cross-sectional averages of cross-products of residuals for two different dates, are useful to check for omitted factors. ${ }^{8}$

We summarize the statistics to test null hypotheses on the number of latent factors next.

Definition 1 The statistics to test the null hypothesis $H_{0}(k)$ of $k$ latent factors are: (a) the squared norm statistic $\mathscr{T}(k):=n \sum_{j=k+1}^{T} \hat{\gamma}_{j}^{2}=n\|\hat{S}\|^{2}$, and $(b)$ the LR statistic $L R(k):=-n \sum_{j=k+1}^{T} \log (1+$ $\left.\hat{\gamma}_{j}\right)=\frac{n}{2}\|\hat{S}\|^{2}+o_{p}(1)$, where $\hat{\gamma}_{k+j}=\delta_{k+j}\left(\hat{V}_{y} \hat{V}_{\varepsilon}^{-1}\right)-1=\delta_{j}(\hat{S})$, for $j=1, \ldots, T-k$, and we denote by $\delta_{j}(\cdot)$ the $j$ th largest eigenvalue of a symmetric matrix.

The statistics in Definition 1 only use the information contained in the eigenvalues of matrix $\hat{S}$. Next we establish the asymptotic distributions of those test statistics with $n \rightarrow \infty$ and $T$ fixed.

## 3 Asymptotic distributional theory

We start by defining the normalization for the latent factor matrix $F=\left[F_{1}: \cdots: F_{k}\right]$ in population. Following classical FA, we set $\mu_{\beta}=0, V_{\beta}=I_{k}$, and $F^{\prime} V_{\varepsilon}^{-1} F=\operatorname{diag}\left(\gamma_{1}, \ldots, \gamma_{k}\right)$, where $V_{\beta}=$ $\lim _{n \rightarrow \infty} \frac{1}{n} \beta^{\prime} \beta$ and $\mu_{\beta}=\lim _{n \rightarrow \infty} \bar{\beta}$ with $\bar{\beta}=\frac{1}{n} \sum_{i=1}^{n} \beta_{i}$. Then, under our assumptions, we have $V_{y}:=$ $\operatorname{plim}_{n \rightarrow \infty} \hat{V}_{y}=F F^{\prime}+V_{\varepsilon}$ and $V_{y} V_{\varepsilon}^{-1} F_{j}=\left(1+\gamma_{j}\right) F_{j}$, i.e., the $F_{j}$ are eigenvectors of matrix $V_{y} V_{\varepsilon}^{-1}$

[^5]associated with eigenvalues $1+\gamma_{j}, j=1, \ldots, k$. ${ }^{9}$ For any given $n$, we define $\tilde{V}_{\varepsilon}=\frac{1}{n} E\left[\varepsilon \varepsilon^{\prime}\right]$ and $\tilde{V}_{\beta}=\frac{1}{n} \beta^{\prime} \beta$, and use a factor normalization in sample that is analogue to the one in population, i.e., $\bar{\beta}=0, \tilde{V}_{\beta}=I_{k}$ (see Assumption A.1) and $F^{\prime} \tilde{V}_{\varepsilon}^{-1} F$ is diagonal. Thus, the normalization of the factor values $F=F_{(n)}$ is sample dependent; we skip index $n$ for the purpose of easing notation.

We use a block-dependence structure to allow for weak cross-sectional dependence in errors.
Assumption 2 (a) The errors are such that $\varepsilon=V_{\varepsilon}^{1 / 2} W \Sigma^{1 / 2}$, where $W=\left[w_{1}: \cdots: w_{n}\right]$ is a $T \times n$ random matrix of standardized errors terms $w_{i, t}$ that are independent across $i$ and uncorrelated across $t$, and $\Sigma=\left(\sigma_{i, j}\right)$ is a positive-definite symmetric $n \times n$ matrix, such that $\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{i=1}^{n} \sigma_{i i}=$ 1. (b) Matrix $\Sigma$ is block diagonal with $J_{n}$ blocks of size $b_{m, n}=B_{m, n} n$, for $m=1, \ldots, J_{n}$, where $J_{n} \rightarrow \infty$ as $n \rightarrow \infty$, and $I_{m}$ denotes the set of indices in block $m$. (c) There exist constants $\delta \in[0,1]$ and $C>0$ such that $\max _{i \in I_{m}} \sum_{j \in I_{m}}\left|\sigma_{i, j}\right| \leq C b_{m, n}^{\delta}$. (d) The block sizes $b_{m, n}$ and block number $J_{n}$ are such that $n^{2 \delta} \sum_{m=1}^{J_{n}} B_{m, n}^{2(1+\delta)}=o(1)$.

As already remarked, the diagonal elements of $V_{\varepsilon}$ are the sample realizations of the common component driving the variance of the error terms at times $t=1, \ldots, T$; see e.g. Barigozzi and Hallin (2016), Renault, Van Der Heijden and Werker (2022) for theory and empirical evidence pointing to variance factors. A sphericity assumption cannot accommodate such a common time-varying component. In empirical applications on individual stocks, blocks in $\Sigma$ can match industrial sectors (Gagliardini, Ossola, and Scaillet (2016)). Assumption 2 (a) is coherent with Assumption 1. Indeed, $\frac{1}{n} E\left[\varepsilon \varepsilon^{\prime}\right]=V_{\varepsilon}^{1 / 2} \frac{1}{n} \sum_{i, j=1}^{n} \sigma_{i, j} E\left[w_{i} w_{j}^{\prime}\right] V_{\varepsilon}^{1 / 2}=\frac{1}{n} \sum_{i=1}^{n} \sigma_{i i} V_{\varepsilon}$ is diagonal, and converges to matrix $V_{\varepsilon}$ in the limit $n \rightarrow \infty$ under the normalization $\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{i=1}^{n} \sigma_{i i}=1$. That normalization is without loss of generality by rescaling of the parameters. Assumption 2 (c) builds on Bickel and Levina (2008), and $\delta<1$ holds under sparsity, vanishing correlations or mixing dependence within blocks. With blocks of equal size, Assumption 2 (d) holds for $J_{n}=n^{\bar{\alpha}}$ and $\bar{\alpha}>\frac{2 \delta}{2 \delta+1}$. Hav-

[^6]ing $\delta<1$ helps relaxing this condition on block granularity, however it is not strictly necessary because we allow value $\delta=1$.

### 3.1 Asymptotic expansions of estimators $\hat{V}_{\varepsilon}$ and $\hat{F}$

Using Equation (1) and the factor normalization in sample, we have $\hat{V}_{y}=\tilde{V}_{y}+\frac{1}{\sqrt{n}} \Psi_{y}+o_{p}\left(\frac{1}{\sqrt{n}}\right)$, where $\tilde{V}_{y}=F F^{\prime}+\tilde{V}_{\varepsilon}$ and $\Psi_{y}=\frac{1}{\sqrt{n}}\left(\varepsilon \beta F^{\prime}+F \beta^{\prime} \varepsilon^{\prime}\right)+\sqrt{n}\left(\frac{1}{n} \varepsilon \varepsilon^{\prime}-\tilde{V}_{\varepsilon}\right)$ (see Appendix D.2). The FA estimators $\hat{V}_{\varepsilon}$ and $\hat{F}$ are consistent M-estimators under nonlinear constraints, and admit expansions at first order for fixed $T$ and $n \rightarrow \infty$, namely $\hat{V}_{\varepsilon}=\tilde{V}_{\varepsilon}+\frac{1}{\sqrt{n}} \Psi_{\varepsilon}+o_{p}\left(\frac{1}{\sqrt{n}}\right)$ and $\hat{F}_{j}=$ $F_{j}+\frac{1}{\sqrt{n}} \Psi_{F_{j}}+o_{p}\left(\frac{1}{\sqrt{n}}\right)$ (see Appendix D.4.1). The next proposition characterizes the diagonal random matrix $\Psi_{\varepsilon}$ and the random vectors $\Psi_{F_{j}}$ by using conditions (FA1) and (FA2) above.

Assumption 3 The non-zero eigenvalues of $V_{y} V_{\varepsilon}^{-1}-I_{T}$ are distinct, i.e., $\gamma_{1}>\ldots>\gamma_{k}>0$.
Proposition 2 Under Assumptions 1-3 and A.1-A.4, we have (a) for $j=1, \ldots, k$

$$
\begin{equation*}
\Psi_{F_{j}}=R_{j}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F_{j}+\Lambda_{j} \Psi_{\varepsilon} V_{\varepsilon}^{-1} F_{j}, \tag{4}
\end{equation*}
$$

where $R_{j}:=\frac{1}{2 \gamma_{j}} P_{F_{j}, V_{\varepsilon}}+\frac{1}{\gamma_{j}} M_{F, V_{\varepsilon}}+\sum_{\ell=1, \ell \neq j}^{k} \frac{1}{\gamma_{j}-\gamma_{\ell}} P_{F_{\ell}, V_{\varepsilon}}$ and $\Lambda_{j}:=-\sum_{\ell=1, \ell \neq j}^{k} \frac{\gamma_{l}}{\gamma_{j}-\gamma_{\ell}} P_{F_{\ell}, V_{\varepsilon}}$ and $P_{F_{j}, V_{\varepsilon}}=F_{j}\left(F_{j}^{\prime} V_{\varepsilon}^{-1} F_{j}\right)^{-1} F_{j}^{\prime} V_{\varepsilon}^{-1}=\frac{1}{\gamma_{j}} F_{j} F_{j}^{\prime} V_{\varepsilon}^{-1}$ is the GLS orthogonal projection onto $F_{j}$. Further, (b) the diagonal matrix $\Psi_{\varepsilon}$ is such that:

$$
\begin{equation*}
\operatorname{diag}\left(M_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) M_{F, V_{\varepsilon}}^{\prime}\right)=0 \tag{5}
\end{equation*}
$$

Equation (4) yields the asymptotic expansion of the eigenvectors by accounting for estimation errors of matrix $\hat{V}_{y} \hat{V}_{\varepsilon}^{-1}$ (first term) and of the normalization constraint (second term). To interpret Equation (5), we can observe that the matrix $M_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) M_{F, V_{\varepsilon}}^{\prime}$ yields the first-order term in the asymptotic expansion of the test statistic $\sqrt{n} \hat{S}$ (up to the left- and right-multiplication by diagonal matrix $V_{\varepsilon}^{-1 / 2}$ ); see Equation (8). Thus, Equation (5) is implied by the property that the diagonal terms of matrix $\hat{S}$ are equal to zero as stated in Proposition 1 (c).

Let us now give the explicit expression of $\Psi_{\varepsilon}$. By using $M_{F, V_{\varepsilon}} \Psi_{y} M_{F, V_{\varepsilon}}^{\prime}=M_{F, V_{\varepsilon}} Z_{n} M_{F, V_{\varepsilon}}^{\prime}$, where $Z_{n}:=\sqrt{n}\left(\frac{1}{n} \varepsilon \varepsilon^{\prime}-\frac{1}{n} E\left[\varepsilon \varepsilon^{\prime}\right]\right)$ is the standardized, centered sample mean of cross-moments of errors, we can rewrite Equation (5) as $\operatorname{diag}\left(M_{F, V_{\varepsilon}}\left(Z_{n}-\Psi_{\varepsilon}\right) M_{F, V_{\varepsilon}}^{\prime}\right)=0$. Now, because $\Psi_{\varepsilon}$ is diagonal, we have $\operatorname{diag}\left(M_{F, V_{\varepsilon}} \Psi_{\varepsilon} M_{F, V_{\varepsilon}}^{\prime}\right)=M_{F, V_{\varepsilon}}^{\odot 2} \operatorname{diag}\left(\Psi_{\varepsilon}\right)$, where $M_{F, V_{\varepsilon}}^{\odot}=M_{F, V_{\varepsilon}}^{\odot} \odot M_{F, V_{\varepsilon}}$ and $\odot$ denotes the Hadamard product (i.e., element-wise matrix product). Thus, we get the equation:

$$
\begin{equation*}
M_{F, V_{\varepsilon}}^{\odot 2} \operatorname{diag}\left(\Psi_{\varepsilon}\right)=\operatorname{diag}\left(M_{F, V_{\varepsilon}} Z_{n} M_{F, V_{\varepsilon}}^{\prime}\right) \tag{6}
\end{equation*}
$$

To have a unique solution for vector $\operatorname{diag}\left(\Psi_{\varepsilon}\right)$, we need the non-singularity of the $T \times T$ matrix on the LHS of this linear equation system. It is the local identification condition in the FA model (see Lemma 7 in Appendix D. 3 i), where we show equivalence with invertibility of the bordered Hessian, i.e., the Hessian of the Lagrangian function in a constrained M-estimation).

Assumption 4 Matrix $M_{F, V_{\varepsilon}}^{\odot 2}$ is non-singular.
Under Assumption 4, we get from Equation (6):

$$
\begin{equation*}
\Psi_{\varepsilon}=\mathcal{T}_{F, V_{\varepsilon}}\left(Z_{n}\right) \tag{7}
\end{equation*}
$$

where $\mathcal{T}_{F, V_{\varepsilon}}(V):=\operatorname{diag}\left(\left[M_{F, V_{\varepsilon}}^{\odot 2}\right]^{-1} \operatorname{diag}\left(M_{F, V_{\varepsilon}} V M_{F, V_{\varepsilon}}^{\prime}\right)\right)$, for any $T \times T$ matrix $V$. Mapping $\mathcal{T}_{F, V_{\varepsilon}}(\cdot)$ is linear and such that $\mathcal{T}_{F, V_{\varepsilon}}(V)=V$, for a diagonal matrix $V$.

Anderson and Rubin (1956), Theorem 12.1, show that the FA estimator is asymptotically normal if $\sqrt{n}\left(\hat{V}_{y}-V_{y}\right)$ is asymptotically normal. They use a linearization of the first-order conditions similar as the one of Proposition 2. Their Equation (12.16) corresponds to our Equation (5). However, they only provide an implicit characterization of the $\Psi_{F_{j}}$ and not an explicit expression for $\Psi_{\varepsilon}$ and $\Psi_{F_{j}}$ in terms of asymptotically Gaussian random matrices like $Z_{n}$ as we do. These key developments pave the way to establishing the asymptotic distributions of estimators $\hat{F}$ and $\hat{V}_{\varepsilon}$ in general settings, that we cover in OA Section D.4, and of the test statistics for the number of factors, that we address next.

### 3.2 Asymptotic expansions of the test statistics

By expanding the terms in the definition of $\hat{S}$ in Equation (2), and using Equation (7) and the $\sqrt{n}$-consistency of FA estimators (see Appendix D.4.1), we have:

$$
\begin{align*}
\sqrt{n} \hat{S} & =V_{\varepsilon}^{-1 / 2} M_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) M_{F, V_{\varepsilon}}^{\prime} V_{\varepsilon}^{-1 / 2}+o_{p}(1) \\
& =V_{\varepsilon}^{-1 / 2} M_{F, V_{\varepsilon}}\left(Z_{n}-\mathcal{T}_{F, V_{\varepsilon}}\left(Z_{n}\right)\right) M_{F, V_{\varepsilon}}^{\prime} V_{\varepsilon}^{-1 / 2}+o_{p}(1) \tag{8}
\end{align*}
$$

Let us now rework the RHS. First, we use that $Z_{n}=\sqrt{n}\left(\frac{1}{n} \varepsilon \varepsilon^{\prime}-\tilde{V}_{\varepsilon}\right)$, where $\tilde{V}_{\varepsilon}=\frac{1}{n} E\left[\varepsilon \varepsilon^{\prime}\right]$ is diagonal and such that $\mathcal{T}_{F, V_{\varepsilon}}\left(\tilde{V}_{\varepsilon}\right)=\tilde{V}_{\varepsilon}$. Thus, we have $Z_{n}-\mathcal{T}_{F, V_{\varepsilon}}\left(Z_{n}\right)=\sqrt{n}\left(\frac{1}{n} \varepsilon \varepsilon^{\prime}-\mathcal{T}_{F, V_{\varepsilon}}\left(\frac{1}{n} \varepsilon \varepsilon^{\prime}\right)\right)=$ : $\bar{Z}_{n}$. We get that $E\left[\bar{Z}_{n}\right]=0$ because diagonal matrices are invariant under mapping $\mathcal{T}_{F, V_{\varepsilon}}(\cdot)$. Second, we write the orthogonal projection as $M_{F, V_{\varepsilon}}=G G^{\prime} V_{\varepsilon}^{-1}$, where $G$ be a $T \times(T-k)$ matrix such that $F^{\prime} V_{\varepsilon}^{-1} G=0$ and $G^{\prime} V_{\varepsilon}^{-1} G=I_{T-k}$. Matrix $G$ is unique up to post-multiplication by an orthogonal matrix. Then, from Equation (8), we get the asymptotic expansion of matrix $\hat{S}$ as

$$
\begin{equation*}
\sqrt{n} \hat{S}=V_{\varepsilon}^{-1 / 2} G \bar{Z}_{n}^{*} G^{\prime} V_{\varepsilon}^{-1 / 2}+o_{p}(1) \tag{9}
\end{equation*}
$$

where $\bar{Z}_{n}^{*}:=G^{\prime} V_{\varepsilon}^{-1} \bar{Z}_{n} V_{\varepsilon}^{-1} G$. The asymptotic distribution of $\sqrt{n} \hat{S}$ is driven by the symmetric $(T-k) \times(T-k)$ zero-mean matrix $\bar{Z}_{n}^{*}$. We have $\operatorname{diag}\left(V_{\varepsilon}^{-1 / 2} G \bar{Z}_{n}^{*} G^{\prime} V_{\varepsilon}^{-1 / 2}\right)=0$ as a consequence of Equation (5) and results above (see also Proposition 1 (c)). We can rewrite the number of degrees of freedom as $d f=\frac{1}{2}(T-k)(T-k+1)-T$, i.e., the number of different elements in $\bar{Z}_{n}^{*}$ minus the number of linear constraints in $\operatorname{diag}\left(V_{\varepsilon}^{-1 / 2} G \bar{Z}_{n}^{*} G^{\prime} V_{\varepsilon}^{-1 / 2}\right)=0$. Matrices $V_{\varepsilon}^{-1 / 2} G \bar{Z}_{n}^{*} G^{\prime} V_{\varepsilon}^{-1 / 2}$ and $\bar{Z}_{n}^{*}$ have the same Frobenius norm because the columns of $V_{\varepsilon}^{-1 / 2} G$ are orthonormal. From Equation (9), the asymptotic expansions for the test statistics in Definition 1 are:

$$
\begin{equation*}
\mathscr{T}(k)=\left\|\bar{Z}_{n}^{*}\right\|^{2}+o_{p}(1), \quad L R(k)=\frac{1}{2}\left\|\bar{Z}_{n}^{*}\right\|^{2}+o_{p}(1) \tag{10}
\end{equation*}
$$

under the null hypothesis of $k$ latent factors. ${ }^{10}$

[^7]We can get further insights in the above results by using the next proposition. Let the $T$ dimensional vectors $g_{j}$ for $j=1, \ldots, T-k$ be the columns of matrix $G$, and let us define the $p \times T$ matrix $\boldsymbol{X}=\left[g_{1} \odot g_{1}: \cdots: g_{T-k} \odot g_{T-k}:\left\{\sqrt{2}\left(g_{i} \odot g_{j}\right)\right\}_{i<j}\right]^{\prime}$, where the pairs of indices $(i, j)$ with $i<j$ are ranked as $(1,2),(1,3), \ldots,(1, T-k),(2,3), \ldots,(T-k-1, T-k)$, and $p=\frac{1}{2}(T-k)(T-k+1)$. Moreover, for a $(T-k) \times(T-k)$ symmetric matrix $Z=\left(z_{i, j}\right)$, let us define the $p$-dimensional vector $\operatorname{vech}(Z)=\left(\frac{1}{\sqrt{2}} z_{11}, \ldots, \frac{1}{\sqrt{2}} z_{T-k, T-k},\left\{z_{i, j}\right\}_{i<j}\right)^{\prime}$, where the out-of-diagonal elements with indices $i<j$ are ranked as above. ${ }^{11}$

Proposition 3 Under Assumptions 1-4, we have (a) $M_{F, V_{\varepsilon}}^{\odot 2}=\boldsymbol{X}^{\prime} \boldsymbol{X} V_{\varepsilon}^{-2}$, (b) $\operatorname{diag}\left(\Psi_{\varepsilon}\right)=$ $\sqrt{2} V_{\varepsilon}^{2}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime} \operatorname{vech}\left(Z_{n}^{*}\right)$, where $Z_{n}^{*}=G^{\prime} V_{\varepsilon}^{-1} Z_{n} V_{\varepsilon}^{-1} G$, and (c) $\operatorname{vech}\left(\bar{Z}_{n}^{*}\right)=$ $\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right) \operatorname{vech}\left(Z_{n}^{*}\right)$.

From Proposition 3 (a), we can state the local identification condition in Assumption 4 as a fullrank condition for matrix $\boldsymbol{X}$ analogously as in linear regression. In part (b), we write the diagonal of $\Psi_{\varepsilon}$ via the coefficients of a OLS regression of the half-vectorization of $Z_{n}^{*}$ onto $\boldsymbol{X}$. Part (c) shows that, after half-vectorization, we can represent the elements of matrix variate $\bar{Z}_{n}^{*}$ as the residual of the orthogonal projection of vech $\left(Z_{n}^{*}\right)$ onto the columns of $\boldsymbol{X}$. Matrix $I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}$ is idempotent of rank $p-T=d f$. Using $\frac{1}{2}\left\|\bar{Z}_{n}^{*}\right\|^{2}=\operatorname{vech}\left(\bar{Z}_{n}^{*}\right)^{\prime} \operatorname{vech}\left(\bar{Z}_{n}^{*}\right)$ and Proposition 3 (c), the leading term in the asymptotic expansions (10) is $\frac{1}{2}\left\|\bar{Z}_{n}^{*}\right\|^{2}=\operatorname{vech}\left(Z_{n}^{*}\right)^{\prime}\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right) \operatorname{vech}\left(Z_{n}^{*}\right)$.

### 3.3 Feasible Central Limit Theorem

We now establish the distributional convergence $\bar{Z}_{n}^{*} \Rightarrow \bar{Z}^{*}$ as $n \rightarrow \infty$ and $T$ is fixed, where $\bar{Z}^{*}$ is a Gaussian symmetric matrix variate. We have that $\bar{Z}^{*}=G^{\prime} V_{\varepsilon}^{-1}\left(Z-\mathcal{T}_{F, V_{\varepsilon}}(Z)\right) V_{\varepsilon}^{-1} G$, where $Z$ is the distributional limit of $Z_{n}$. Establishing a feasible Central Limit Theorem (CLT) via a nonparametric estimator of the asymptotic variance is easier for $\bar{Z}_{n}^{*}$ than for $Z_{n}$, and is sufficient for

[^8]testing purposes. ${ }^{12}$ By the block structure in Assumption 2 (b), we can write $\bar{Z}_{n}^{*}$ as a sum of independent zero-mean terms: $\bar{Z}_{n}^{*}=\frac{1}{\sqrt{n}} G^{\prime} V_{\varepsilon}^{-1}\left(\varepsilon \varepsilon^{\prime}-\mathcal{T}_{F, V_{\varepsilon}}\left(\varepsilon \varepsilon^{\prime}\right)\right) V_{\varepsilon}^{-1} G=\frac{1}{\sqrt{n}} \sum_{m=1}^{J_{n}} z_{m, n}$, where the variables in the triangular array $z_{m, n}=\sum_{i, j \in I_{m}} \sigma_{i, j} G^{\prime} V_{\varepsilon}^{-1 / 2}\left[w_{i} w_{j}^{\prime}-\mathcal{T}_{F, V_{\varepsilon}}\left(w_{i} w_{j}^{\prime}\right)\right] V_{\varepsilon}^{-1 / 2} G=$ $\sum_{i \in I_{m}} G^{\prime} V_{\varepsilon}^{-1}\left[\varepsilon_{i} \varepsilon_{i}^{\prime}-\mathcal{T}_{F, V_{\varepsilon}}\left(\varepsilon_{i} \varepsilon_{i}^{\prime}\right)\right] V_{\varepsilon}^{-1} G$ are independent across $m$ and such that $E\left[z_{m, n}\right]=0$. In Appendix B, we invoque the CLT for independent heterogeneous variables to vech $\left(\bar{Z}_{n}^{*}\right)=$ $\frac{1}{\sqrt{n}} \sum_{m=1}^{J_{n}} \operatorname{vech}\left(z_{m, n}\right)$ and use Assumptions 2 (c) and (d) to check the Liapunov condition. Then, we get $\bar{Z}_{n}^{*} \Rightarrow \bar{Z}^{*}$, where $\operatorname{vech}\left(\bar{Z}^{*}\right) \sim N\left(0, \Omega_{\bar{Z}^{*}}\right)$ and $\Omega_{\bar{Z}^{*}}=\lim _{n \rightarrow \infty^{n}} \frac{1}{n} \sum_{m=1}^{J_{n}} V\left[\operatorname{vech}\left(z_{m, n}\right)\right]$. Our assumptions imply that $\Omega_{\bar{Z}^{*}}$ is finite. From Proposition 3 (c), we have $\Omega_{\bar{Z}^{*}}=\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right)$ $\Omega_{Z^{*}}\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right)$, where $\operatorname{vech}\left(Z^{*}\right) \sim N\left(0, \Omega_{Z^{*}}\right)$, for $Z^{*}=G^{\prime} V_{\varepsilon}^{-1} Z V_{\varepsilon}^{-1} G$. We characterize the variance $\Omega_{Z}=V[\operatorname{vech}(Z)]$ of the distributional limit of $Z_{n}$ in Lemma 1 in Appendix B. In particular, matrix $\Omega_{\bar{Z}^{*}}$ is singular with rank $d f$. Then, the asymptotic expansions in (10) yield the asymptotic distributions for the test statistics.

Proposition 4 Let Assumptions 1-4 and A.1-A. 5 hold. As $n \rightarrow \infty$ and $T$ is fixed, under the null hypothesis $H_{0}(k)$ of $k$ latent factors, (a) $\mathscr{T}(k) \Rightarrow\left\|\bar{Z}^{*}\right\|^{2}, L R(k) \Rightarrow \frac{1}{2}\left\|\bar{Z}^{*}\right\|^{2}$, where vech $\left(\bar{Z}^{*}\right)$ $\sim N\left(0, \Omega_{\bar{Z}^{*}}\right)$, and (b) $\hat{\mathscr{R}} \hat{\Omega}_{\bar{Z}^{*}} \hat{\mathscr{R}}^{-1} \xrightarrow{p} \Omega_{\bar{Z}^{*}}$, where $\hat{\Omega}_{\bar{Z}^{*}}=\frac{1}{n} \sum_{m=1}^{J_{n}} \operatorname{vech}\left(\hat{z}_{m, n}\right)$ vech $\left(\hat{z}_{m, n}\right)^{\prime}$ and $\hat{z}_{m, n}=\sum_{i \in I_{m}} \hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1}\left(\hat{\varepsilon}_{i} \hat{\varepsilon}_{i}^{\prime}-\mathcal{T}_{\hat{F}, \hat{V}_{\varepsilon}}\left(\hat{\varepsilon}_{i} \hat{\varepsilon}_{i}^{\prime}\right)\right) \hat{V}_{\varepsilon}^{-1} \hat{G}$, with $\hat{\varepsilon}_{i}=M_{\hat{F}, \hat{V}_{\varepsilon}}\left(y_{i}-\bar{y}\right)$, for an orthogonal matrix $\hat{\mathscr{R}}$. Under the alternative hypothesis $H_{1}(k)$ of more than $k$ latent factors, (c) $\mathscr{T}(k) \geq C n$ and $L R(k) \geq C n$, w.p.a. 1 for a constant $C>0$, and $\hat{\Omega}_{\bar{Z}^{*}}=O_{p}\left(n \sum_{m=1}^{J_{n}} B_{m, n}^{2}\right)=o_{p}(n)$.

From Proposition 4 (a), using $\frac{1}{2}\left\|\bar{Z}^{*}\right\|^{2}=\operatorname{vech}\left(\bar{Z}^{*}\right)^{\prime} \operatorname{vech}\left(\bar{Z}^{*}\right) \sim \sum_{j=1}^{d f} \mu_{j} \chi_{j}^{2}(1)$, the asymptotic distribution of the LR statistic is a weighted average of $d f$ mutually independent chi-square variates with weights $\mu_{j}$ that are the non-zero eigenvalues of matrix $\Omega_{\bar{Z}^{*}} .{ }^{13}$ In Proposition 4 (b), we use $\hat{G}=\hat{V}_{\varepsilon}^{1 / 2} \hat{Q}$, where $\hat{Q}$ is a $T \times(T-k)$ matrix with orthonormal columns that span the range

[^9]of $I_{T}-\hat{V}_{\varepsilon}^{-1 / 2} \hat{F}\left(\hat{F}^{\prime} \hat{V}_{\varepsilon}^{-1} \hat{F}\right)^{-1} \hat{F}^{\prime} \hat{V}_{\varepsilon}^{-1 / 2}$. ${ }^{14}$ The orthogonal matrix $\hat{\mathscr{R}}$ accounts for an arbitrary choice of that orthonormal basis. With fixed $T$, the GLS residuals $\hat{\varepsilon}_{i}$ are asymptotically close to $M_{F, V_{\varepsilon}} \varepsilon_{i}$ and not to the true errors $\varepsilon_{i}$. This fact does not impede the consistency of $\hat{\Omega}_{\bar{Z}^{*}}$, because $G^{\prime} V_{\varepsilon}^{-1} M_{F, V_{\varepsilon}}=G^{\prime} V_{\varepsilon}^{-1}$ and $\mathcal{T}_{F, V_{\varepsilon}}\left(M_{F, V_{\varepsilon}} V M_{F, V_{\varepsilon}}^{\prime}\right)=\mathcal{T}_{F, V_{\varepsilon}}(V)$, for any $V$. We can consistently estimate the critical values of the asymptotic statistics $\left\|\bar{Z}^{*}\right\|^{2}$ by simulating a large number of draws from a Gaussian symmetric matrix variate with vectorized variance $\hat{\Omega}_{\bar{Z}^{*}}$, whose norms are unaffected by the orthogonal matrix $\hat{\mathscr{R}}$. Finally, Proposition 4 (c) gives test consistency against global alternative hypotheses.

## 4 Discussion of three special cases

In this section, we particularize the general distributional results of Proposition 4 to three important cases, namely Gaussian errors, settings where the asymptotic distribution under Gaussian errors still holds for the test statistics (up to scaling), and spherical errors.

### 4.1 Gaussian errors

Let us consider the case where the errors $\varepsilon_{i} \stackrel{i n d}{\sim} N\left(0, \sigma_{i i} V_{\varepsilon}\right)$ are independent Gaussian vectors. From classical FA theory, we expect that the statistic $L R(k)$ admits asymptotically a chi-square distribution with $d f$ degrees of freedom in the cross-sectionally homoschedastic case, i.e., $\sigma_{i i}=1$ for all assets $i$. We cannot expect that this distributional result applies to the Gaussian framework in full generality, since - even in such a case - our setting corresponds to a pseudo model (because the $\sigma_{i i}$ may be heterogeneous across $i$, and the $\beta_{i}$ are treated as fixed effects, namely incidental parameters, instead of Gaussian random effects). In order to establish the distribution of $\frac{1}{2}\left\|\bar{Z}^{*}\right\|^{2}=$ $\operatorname{vech}\left(\bar{Z}^{*}\right)^{\prime} \operatorname{vech}\left(\bar{Z}^{*}\right)$, we use Proposition 3 (c) written for the distributional limits to get $\frac{1}{2}\left\|\bar{Z}^{*}\right\|^{2}=$

[^10]vech $\left(Z^{*}\right)^{\prime}\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right)$ vech $\left(Z^{*}\right)$, where $I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}$ is idempotent of rank $d f$. Under the normality assumption for the error terms, we have $\varepsilon_{i}^{*}:=G^{\prime} V_{\varepsilon}^{-1} \varepsilon_{i} \stackrel{i n d}{\sim} N\left(0, \sigma_{i i} I_{T-k}\right)$. Thus, by the Liapunov CLT, the distributional limit of $\frac{1}{\sqrt{q}} Z_{n}^{*}=\sqrt{n / q}\left(\frac{1}{n} \varepsilon^{*}\left(\varepsilon^{*}\right)^{\prime}-E\left[\frac{1}{n} \varepsilon^{*}\left(\varepsilon^{*}\right)^{\prime}\right]\right)$ is in the Gaussian Orthogonal Ensemble (GOE) for dimension $T-k$ (see e.g. Tao (2012)), i.e., $\frac{1}{\sqrt{q}} \operatorname{vech}\left(Z^{*}\right) \sim N\left(0, I_{p}\right)$, where $q:=\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{i=1}^{n} \sigma_{i i}^{2}$. Then, we get $L R(k) \Rightarrow \frac{1}{2}\left\|\bar{Z}^{*}\right\|^{2} \sim q \chi^{2}(d f)$, i.e., a scaled chi-square distribution with $d f$ degrees of freedom. In the cross-sectionally homoschedastic case, we have $q=1$ yielding the classical $\chi^{2}(d f)$ result. On the contrary, crosssectional heterogeneity in the unconditional idiosyncratic variances yields $q>1$ and a deviation from classical FA theory even in the Gaussian case. Hence, unobserved heterogeneity across asset idiosyncratic variances would lead to an oversized LR test if we use critical values from the chi-square table without proper scaling.

### 4.2 Validity of the scaled asymptotic chi-square test

In this subsection, we investigate sufficient conditions for the validity of the scaled asymptotic $\chi^{2}(d f)$ distribution of the LR statistic in special cases beyond Gaussianity of errors. For this purpose, let us first notice that $\bar{Z}=Z-\mathcal{T}_{F, V_{\varepsilon}}(Z)$ only involves the out-of-diagonal elements of $Z$. Under independent Gaussian errors (Section 4.1), by the Liapunov CLT, we have $Z_{t, s} \sim$ $N\left(0, q V_{\varepsilon, t t} V_{\varepsilon, s s}\right)$, for $t>s$, mutually independent, where the $V_{\varepsilon, t t}$ are the diagonal elements of matrix $V_{\varepsilon}$. We deduce that any setting featuring the same joint asymptotic distribution for the out-of-diagonal elements of random matrix $Z_{n}$ leads to the same asymptotic distribution of the LR statistic as in the Gaussian case, namely the scaled $\chi^{2}(d f)$ distribution.

Proposition 5 Let Assumptions 1-4, A.1-A.4 hold with (a) $\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{i=1}^{n} E\left[\varepsilon_{i, t} \varepsilon_{i, s} \varepsilon_{i, r} \varepsilon_{i, p}\right]=q V_{\varepsilon, t} V_{\varepsilon, s s}$, when $t=r>s=p$, for a constant $q>0$, and $=0$ in all other cases with $t>s$ and $r>p$, and (b) let $\kappa=\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{m=1}^{J_{n}} \sum_{i \neq j \in I_{m}} \sigma_{i j}^{2}$ as in Assumption A.4 (b). Then, $L R(k) \Rightarrow \bar{q} \chi^{2}(d f)$ under $H_{0}(k)$ for $\bar{q}:=q+\kappa$.

Conditions (a) and (b) in Proposition 5 generalize the correctness of the scaled chi-square test beyond Gaussianity and error independence across time and assets. Under Assumption 2, Condition (a) is satisfied if the standardized error terms $w_{i, t}$ are conditionally homoschedastic martingale difference sequences. However, Condition (a) excludes empirically relevant cases such as ARCH processes for $w_{i, t}$, because, in that case, $\frac{1}{V_{\varepsilon, t t} V_{\varepsilon, s s}} E\left[\varepsilon_{i, t}^{2} \varepsilon_{i, s}^{2}\right]$ depends on lag $t-s$. Hence, serial correlation in squared idiosyncratic errors is responsible for the deviation of the LR test from the scaled chi-square asymptotic distribution. This setting is covered by the general results in Proposition 4.

Anderson and Amemiya (1988) establish the asymptotic distribution of FA estimates assuming that the error terms are i.i.d. across sample units and deploy an assumption that is analogue to Condition (a) above in their Corollary 2. The i.i.d. assumption in our case implies $\sigma_{i i}=1$ for all $i$, which results in a cross-sectionally homoschedastic setting. ${ }^{15}$ That setting is irrealistic in our application, as it would imply that the idiosyncratic variance is the same for all assets. Our results show that establishing the asymptotic distribution of the test statistics, especially the AUMPI property of LR test (see Section 5), in a general setting with non-Gaussian errors, heterogeneous idiosyncratic variances and ARCH effects, is challenging, but still possible.

### 4.3 Spherical errors

When errors are spherical, i.e., matrix $V_{\varepsilon}=\bar{\sigma}^{2} I_{T}$ is a multiple of the identity with unknown parameter $\bar{\sigma}^{2}>0$, and this restriction on $V_{\varepsilon}$ is imposed in the estimation procedure, the FA estimator $\hat{F}$ boils down to the Principal Component Analysis (PCA) estimator; see Anderson and Rubin (1956) Section 7.3. Then, $\hat{F}$ is the matrix of eigenvectors of matrix $\hat{V}_{y}$ standardized such that $\hat{F}^{\prime} \hat{F}=$ $\operatorname{diag}\left(\hat{\delta}_{1}-\hat{\sigma}^{2}, \ldots, \hat{\delta}_{k}-\hat{\sigma}^{2}\right)$, and $\hat{\sigma}^{2}=\frac{1}{T-k} \sum_{j=k+1}^{T} \hat{\delta}_{j}$, where $\hat{\delta}_{j}=\delta_{j}\left(\hat{V}_{y}\right)$. The statistic $\hat{S}$ becomes

[^11]$\hat{S}=\frac{1}{\hat{\sigma}^{2}} M_{\hat{F}}\left(\hat{V}_{y}-\hat{\sigma}^{2} I_{T}\right) M_{\hat{F}}=\frac{1}{\hat{\sigma}^{2}}\left(\frac{1}{n} \hat{\varepsilon} \hat{\varepsilon}^{\prime}\right)-M_{\hat{F}}$, where $M_{\hat{F}}=I_{T}-\hat{F}\left(\hat{F}^{\prime} \hat{F}\right)^{-1} \hat{F}^{\prime}$ and $\hat{\varepsilon}=M_{\hat{F}} \tilde{Y}$ is the matrix of OLS residuals. Then, the statistic $L R(k)$ boils down to the LR statistic invoqued by Onatski (2022) in his discussion of FGS. By repeating the arguments of Sections 2 and 3 in the constrained setting of spherical errors, we get $\operatorname{Tr}\left(M_{F}\left(\Psi_{y}-\Psi_{\varepsilon}\right) M_{F}\right)=0$ instead of Equation (5). It yields the asymptotic expansion $\hat{V}_{\varepsilon}=\tilde{\sigma}^{2} I_{T}+\frac{1}{\sqrt{n}} \Psi_{\varepsilon}+O_{p}(1 / n)$, where $\tilde{\sigma}^{2}=\bar{\sigma}^{2} \frac{1}{n} \sum_{i=1}^{n} \sigma_{i i}$ and $\Psi_{\varepsilon}=\frac{1}{T-k} \operatorname{Tr}\left(M_{F} Z_{n}\right) I_{T}$. We get the asymptotic distribution $\sqrt{n} \hat{S} \Rightarrow \frac{1}{\bar{\sigma}^{2}} G \bar{Z}^{*} G^{\prime}$ as $n \rightarrow \infty$ and $T$ is fixed, where we have $\bar{Z}^{*}=\frac{1}{\bar{\sigma}^{4}} G^{\prime}\left(Z-\frac{1}{T-k} \operatorname{Tr}\left(M_{F} Z\right) I_{T}\right) G=Z^{*}-\frac{1}{T-k} \operatorname{Tr}\left(Z^{*}\right) I_{T-k}$, and $G$ is a $T \times(T-k)$ matrix such that $F^{\prime} G=0$ and $G^{\prime} G=\bar{\sigma}^{2} I_{T-k}$. It yields the asymptotic distribution of test statistic $L R(k) \Rightarrow \frac{1}{2}\left(\operatorname{Tr}\left[\left(Z^{*}\right)^{2}\right]-\frac{1}{T-k}\left[\operatorname{Tr}\left(Z^{*}\right)\right]^{2}\right)$ obtained by Onatski (2022). ${ }^{16}$

## 5 Local asymptotic power

In this section, we study the asymptotic power of the test statistics against local alternative hypotheses in which we have $k$ (strong) factors plus a weak factor. Specifically, under $H_{1, l o c}(k)$, we have $\sqrt{n} \gamma_{k+1} \rightarrow c_{k+1}$ as $n \rightarrow \infty$, with $c_{k+1}>0$. The (drifting) DGP is $Y=\mu 1_{n}^{\prime}+F \beta^{\prime}+F_{k+1} \beta_{l o c}^{\prime}+\varepsilon$, where $\beta_{\text {loc }}$ is the loading vector for the $(k+1)$ th factor, and the factor vector is normalized such that $F_{k+1}=\sqrt{\gamma_{k+1}} \rho_{k+1}$ with $\rho_{k+1}^{\prime} V_{\varepsilon}^{-1} \rho_{k+1}=1$ and $F^{\prime} V_{\varepsilon}^{-1} \rho_{k+1}=0$. Thus, we can write $\rho_{k+1}=G \xi_{k+1}$ for a $T-k$ dimensional vector $\xi_{k+1}$ with unit norm. Scalar $c_{k+1}$ and vector $\xi_{k+1}$ yield the (normalized) strength and the direction of the local alternative.

### 5.1 Asymptotic distributions under local alternative hypotheses

The derivation of the asymptotic distribution of $\hat{S}$ under $H_{1, l o c}(k)$ uses the asymptotic expansion $\hat{V}_{y}=\tilde{V}_{y}+\frac{1}{\sqrt{n}} \Psi_{y, l o c}+o_{p}\left(\frac{1}{\sqrt{n}}\right)$, where $\Psi_{y, l o c}=c_{k+1} \rho_{k+1} \rho_{k+1}^{\prime}+\frac{1}{\sqrt{n}}\left(\varepsilon \beta F^{\prime}+F \beta^{\prime} \varepsilon^{\prime}\right)+\sqrt{n}\left(\frac{1}{n} \varepsilon \varepsilon^{\prime}-\tilde{V}_{\varepsilon}\right)$

[^12]and $\tilde{V}_{y}=F F^{\prime}+\tilde{V}_{\varepsilon}$ (see the proof of Proposition 6 in Appendix B for the derivation). Thus, the arguments deployed in Section 3 now apply with $\Psi_{y, l o c}$ instead of $\Psi_{y}$ and lead to the next result.

Proposition 6 Let Assumptions 1-4, A.1-A. 4 hold. Under the local alternative hypothesis $H_{1, l o c}(k)$, we have as $n \rightarrow \infty$ and $T$ is fixed (a) $\sqrt{n} \hat{S} \Rightarrow V_{\varepsilon}^{-1 / 2} G \bar{Z}_{l o c}^{*} G^{\prime} V_{\varepsilon}^{-1 / 2}$, and $(b) \mathscr{T}(k) \Rightarrow\left\|\bar{Z}_{l o c}^{*}\right\|^{2}$ and $L R(k) \Rightarrow \frac{1}{2}\left\|\bar{Z}_{l o c}^{*}\right\|^{2}$, with Gaussian matrix $\bar{Z}_{l o c}^{*}:=\bar{Z}^{*}+\Delta$ and vector vech $(\Delta):=$ $c_{k+1}\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right) \operatorname{vech}\left(\xi_{k+1} \xi_{k+1}^{\prime}\right)$.

Matrix variate $\bar{Z}_{l o c}^{*}$ is a non-central symmetric Gaussian matrix. The non-zero mean depends in general on both $c_{k+1}$ and $\xi_{k+1}$, while the variances and covariances of the elements of $\bar{Z}_{l o c}^{*}$ are the same as those of $\bar{Z}^{*}$. The non-centrality term $\operatorname{vech}(\Delta)$ is in charge of the asymptotic local power of the statistics. When this vector is null, the asymptotic local power is zero. Indeed, for some local "alternatives" the $(k+1)$ th weak factor can be absorbed in the diagonal variance matrix $V_{\varepsilon}$ of the error terms. More precisely, in Appendix D. 3 ii), we show that $V_{y}+\frac{c_{k+1}}{\sqrt{n}} \rho_{k+1} \rho_{k+1}^{\prime}=$ $F^{*}\left(F^{*}\right)^{\prime}+V_{\varepsilon}^{*}+\frac{1}{\sqrt{n}} G \Delta G^{\prime}+o(1 / \sqrt{n})$ for some $T \times k$ matrix $F^{*}$ and diagonal matrix $V_{\varepsilon}^{*}$, which yields asymptotically a $k$-factor model when $\Delta=0$.

Using the expression $\frac{1}{2}\left\|\bar{Z}_{l o c}^{*}\right\|^{2}=\left[\operatorname{vech}\left(Z^{*}\right)+c_{k+1} \operatorname{vech}\left(\xi_{k+1} \xi_{k+1}^{\prime}\right)\right]^{\prime}\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right)$ $\left[\operatorname{vech}\left(Z^{*}\right)+c_{k+1} \operatorname{vech}\left(\xi_{k+1} \xi_{k+1}^{\prime}\right)\right]$ and $\operatorname{vech}\left(Z^{*}\right) \sim N\left(0, \Omega_{Z^{*}}\right)$, from Proposition 6 we deduce that the asymptotic distribution of the $L R(k)$ statistic under the local alternative is a weighted average of $d f$ mutually independent non-central chi-square distributions:

$$
\begin{equation*}
L R(k) \Rightarrow \sum_{j=1}^{d f} \mu_{j} \chi^{2}\left(1, \lambda_{j}^{2}\right) \tag{11}
\end{equation*}
$$

for $\lambda_{j}^{2}=c_{k+1}^{2} \mu_{j}^{-1}\left[v_{j}^{\prime}\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right) \operatorname{vech}\left(\xi_{k+1} \xi_{k+1}^{\prime}\right)\right]^{2}$, where the $\mu_{j}$ and $v_{j}$ are the nonzero eigenvalues and the associated standardized eigenvectors of matrix $\Omega_{\bar{Z}^{*}}$. We have $\lambda^{2}:=$ $\sum_{j=1}^{d f} \mu_{j} \lambda_{j}^{2}=c_{k+1}^{2} \operatorname{vech}\left(\xi_{k+1} \xi_{k+1}^{\prime}\right)^{\prime}\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right) \operatorname{vech}\left(\xi_{k+1} \xi_{k+1}^{\prime}\right)=\frac{1}{2}\|\Delta\|^{2}$, i.e., the half squared Frobenius norm of the matrix measuring local distance from the $k$-factor specification. It follows that the asymptotic local power of the LR statistic is non null as long as $\lambda^{2}>0$, i.e., it has
non-trivial asymptotic power against any proper local alternative hypothesis. In our Monte Carlo experiments reported in the SMC, we find that the LR statistic has size close to the nominal value, and power against global as well as local alternatives with time dimension as small as $T=6$.

Under the normality of errors, or more generally the conditions of Proposition 5, using that matrix $\frac{1}{\sqrt{\bar{q}}} Z^{*}$ is in the GOE for dimension $T-k$, i.e. $\operatorname{vech}\left(Z^{*}\right) \sim N\left(0, \bar{q} I_{p}\right)$, we have $L R(k) \Rightarrow$ $\bar{q} \chi^{2}\left(d f, \lambda^{2} / \bar{q}\right)$ from (11). The local power is a function solely of the squared Euclidean norm of the vector $\operatorname{vech}(\Delta)$ measuring local distance from the $k$-factor specification, divided by $\bar{q}$.

### 5.2 AUMPI tests

In this subsection, we investigate asymptotic local optimality of the LR statistic for testing hypotheses on the number of latent factors. In our framework with composite null and alternative hypotheses and multi-dimensional parameter, we cannot expect in general to establish Uniformly Most Powerful (UMP) tests. Instead, we can establish an optimality property by restricting the class of tests to invariant tests (e.g. Lehmann and Romano (2005)). We focus on statistics with test functions $\phi$ written on the elements of matrix $\hat{S}$. To eliminate the asymptotic redundancy in the elements of $\hat{S}$, we actually consider the test class $\mathscr{C}=\{\phi: \phi=\phi(\hat{W})\}$ with $\hat{W}:=\sqrt{n} \boldsymbol{D}^{\prime} \operatorname{vech}\left(\hat{S}^{*}\right)$, where $\hat{S}^{*}=\hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1 / 2} \hat{S} \hat{V}_{\varepsilon}^{-1 / 2} \hat{G}$ and $\boldsymbol{D}$ is a $p \times d f$ full-rank matrix, such that $I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}=\boldsymbol{D} \boldsymbol{D}^{\prime}$ and $\boldsymbol{D}^{\prime} \boldsymbol{D}=I_{d f}$. Symmetric matrix $\hat{S}^{*}$ contains the information in $\hat{S}=\hat{V}_{\varepsilon}^{-1 / 2} \hat{G} \hat{S}^{*} \hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1 / 2}$ beyond orthogonality to $\hat{V}_{\varepsilon}^{-1 / 2} \hat{F}$. Vector $\hat{W}$ contains the information in $\sqrt{n}$ vech $\left(\hat{S}^{*}\right)$ beyond asymptotic orthogonality to $\boldsymbol{X}$. From Proposition 6 we have $\hat{W} \Rightarrow$ $N\left(0, \boldsymbol{D}^{\prime} \Omega_{Z^{*}} \boldsymbol{D}\right)$ under the null hypothesis $H_{0}(k)$ and $\hat{W} \Rightarrow N\left(c_{k+1} \boldsymbol{D}^{\prime} \operatorname{vech}\left(\xi_{k+1} \xi_{k+1}^{\prime}\right), \boldsymbol{D}^{\prime} \Omega_{Z^{*}} \boldsymbol{D}\right)$ under the local alternative $H_{1, l o c}(k)$, i.e., an asymptotic Gaussian testing problem.

Matrices $G$ and $\boldsymbol{D}$ are both defined up to post-multiplication by an orthogonal matrix. This point yields a group of orthogonal transformations under which we require the test statistics to be invariant. ${ }^{17}$ In Appendix D.5, we show that the maximal invariant under this group is provided

[^13]by $\hat{W}^{\prime} \hat{W}=n \cdot \operatorname{vech}\left(\hat{S}^{*}\right)^{\prime}\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right) \operatorname{vech}\left(\hat{S}^{*}\right)$. Because $\sqrt{n} \cdot \operatorname{vech}\left(\hat{S}^{*}\right)$ belongs to the range of matrix $I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}$ up to $o_{p}(1)$ terms under both $H_{0}(k)$ and $H_{1, l o c}(k)$, we have $\hat{W}^{\prime} \hat{W}=\frac{n}{2}\left\|\hat{S}^{*}\right\|^{2}+o_{p}(1)=\frac{n}{2}\|\hat{S}\|^{2}+o_{p}(1)=\frac{1}{2} \mathscr{T}(k)+o_{p}(1)$. Therefore, the invariant tests are functions of the squared norm statistic $\mathscr{T}(k)$, which is asymptotically equivalent to the LR statistic (up to the factor $1 / 2$ ).

In the Gaussian case, or more generally under the conditions of Proposition 5, the LR statistic follows asymptotically a scaled non-central chi-square distribution with $d f$ degrees of freedom and non-centrality parameter $\lambda^{2}=\sum_{j=1}^{d f} \lambda_{j}^{2}$ as shown in the previous subsection. Thus, we can simplify the null and alternative hypotheses of our testing problem asymptotically and locally to a one-sided test with null hypothesis $H_{0}(k): \lambda^{2}=0$ vs. alternative hypothesis $H_{1, l o c}(k): \lambda^{2}>0$. The scaling constant $q>0$ plays no role in the power analysis. It means that the LR test is an AUMPI test (Lehmann and Romano (2005) Chapters 3 and 13). Indeed, the density $g\left(z ; d f, \lambda^{2}\right)$ of the $\chi^{2}\left(d f, \lambda^{2}\right)$ distribution is Totally Positive of order $2(T P 2)$ in $z$ and $\lambda^{2}$ (Eaton (1987) Example A. 1 p. 468); see Miravete (2011) for a review of applications of TP2 in economics. A density, which is TP2 in $z$ and $\lambda^{2}$, has the Monotone Likelihood Ratio (MLR) property (Eaton (1987) p. 467). Since $g\left(z ; d f, \lambda^{2}\right) / g(z ; d f, 0)$ is an increasing function in $z$, it gives the AUMPI property.

In the general case with $d f>1$, when neither Gaussianity nor the conditions of Proposition 5 apply, we cannot use the same reasoning, since the density $f\left(z ; \lambda_{1}, \ldots, \lambda_{d f}\right)$ of $\sum_{j=1}^{d f} \mu_{j} \chi^{2}\left(1, \lambda_{j}^{2}\right)$, with $\mu_{j}>0, j=1, \ldots, d f$, is not a function of $\lambda^{2}=\sum_{j=1}^{d f} \mu_{j} \lambda_{j}^{2}$ only, and thus cannot be TP2 in $z$ and $\lambda^{2}$. Instead, we use a power series representation of the density of $\sum_{j=1}^{d f} \mu_{j} \chi^{2}\left(1, \lambda_{j}^{2}\right)$ in terms of central chi-square densities from Kotz, Johnson, and Boyd (1967). Under the sufficient condition (12) in Proposition 7, the density ratio $\frac{f\left(z ; \lambda_{1}, \ldots, \lambda_{d f}\right)}{f(z ; 0, \ldots, 0)}$ is monotone increasing in $z$.

Proposition 7 Let Assumptions 1-4, A.1-A.4 hold. (a) Let us assume that, for any DGP in the

[^14]subset $\bar{H}_{1, l o c}(k) \subset H_{1, l o c}(k)$ of the local alternative hypothesis, we have for any integer $m \geq 3$ :
\[

$$
\begin{equation*}
\sum_{j>l \geq 0, j+l=m} \frac{(j-l) \Gamma\left(\frac{d f}{2}\right)^{2}}{\Gamma\left(\frac{d f}{2}+j\right) \Gamma\left(\frac{f f}{2}+l\right)}\left[c_{j}\left(\lambda_{1}, \ldots, \lambda_{d f}\right) c_{l}(0, \ldots, 0)-c_{l}\left(\lambda_{1}, \ldots, \lambda_{d f}\right) c_{j}(0, \ldots, 0)\right] \geq 0 \tag{12}
\end{equation*}
$$

\]

where $\Gamma(\cdot)$ is the Gamma function, $c_{j}\left(\lambda_{1}, \ldots, \lambda_{d f}\right):=E\left[Q\left(\lambda_{1}, \ldots, \lambda_{d f}\right)^{j}\right] / j!$ for $Q\left(\lambda_{1}, \ldots, \lambda_{d f}\right)=$ $\frac{1}{2} \sum_{j=1}^{d f}\left(\sqrt{\nu_{j}} X_{j}+\sqrt{1-\nu_{j}} \lambda_{j}\right)^{2}, \nu_{j}=1-\frac{1}{\mu_{j}} \mu_{1}$ with the $\mu_{j}$ ranked in increasing order, and $X_{j} \sim$ $N(0,1)$ are mutually independent. Then, the statistics $\mathscr{T}(k)$ and $L R(k)$ yield AUMPI tests against $\bar{H}_{1, \text { loc }}(k)$. (b) Suppose that either $\lambda_{1}^{2}+\left(1-\nu_{2}\right) \lambda_{2}^{2} \geq \nu_{2}$ and $\left(1-\nu_{2}\right) \lambda_{2}^{2} \geq \frac{1}{2} \nu_{2}$ when $d f=2$, or

$$
\begin{equation*}
1\{i=0\} \lambda_{1}^{2}+\sum_{j=2}^{d f-1} \rho_{j}^{i}\left(1-\nu_{j}\right) \lambda_{j}^{2}+\left(1-\nu_{d f}\right) \lambda_{d f}^{2} \geq \frac{\nu_{d f}}{i+1}\left(d f-2-\sum_{j=2}^{d f-1} \rho_{j}^{i+1}\right) \tag{13}
\end{equation*}
$$

for all $i \geq 0$, where $\rho_{j}:=\frac{\nu_{j}}{\nu_{d f}}$, when $d f \geq 3$. Then, Inequalities (12) hold for any $m \geq 3$.
Conditions (12) involve polynomial inequalities in the parameters $\lambda_{j}$ of the alternative hypothesis, and parameters $\nu_{j}$ of the weights of the non-central chi-square distributions, $j=1, \ldots, d f$. It is challenging to establish an explicit characterization of the $\lambda_{j}$ and $\nu_{j}$ equivalent to Inequalities (12), unless $d f=1 .^{18}$ By deploying a novel characterization of the $c_{j}\left(\lambda_{1}, \ldots, \lambda_{d f}\right)$ in terms of a recurrence relation (Lemma 3), we establish explicit sufficient conditions in part b) of Proposition 7. Inequalities (13) are linear in the $\lambda_{j}^{2}$, and define a non-empty convex domain in the $\left(\lambda_{1}^{2}, . ., \lambda_{d f}^{2}\right)$ space, that does not contain the origin $\lambda_{1}=\ldots=\lambda_{d f}=0$ (unless the DGP is such that $\nu_{2}=\ldots=\nu_{d f}$, in which case the RHS of (13) is nil for all $i$ and thus any $\lambda_{j}^{2}$ meet the inequalities). Proposition 7 b ) implies that, for a given set of values of $d f$, the MLR property holds if $\lambda_{j} \geq \underline{\lambda}$ for all $j$, uniformly for $\nu_{j} \leq \bar{\nu}$, where $\underline{\lambda}>0$ is a constant that depends on $\bar{\nu}<1$. Vanishing values of the $\nu_{j}$ correspond to homogenous weights $\mu_{j}$, i.e., the scaled non-central chi-square distribution with $d f$ degrees of freedom. Hence, the AUMPI property in Proposition 7 holds in neighborhoods of DGPs that match the conditions of Proposition 5 (e.g. Gaussian errors) for alternative hypotheses that are sufficiently separated from the null hypothesis. Besides, Proposition 7 shows that the

[^15]Gaussian case is not the only design delivering an AUMPI test. Further, in the SMC, we establish an analytical representation of the coefficients $c_{k}\left(\lambda_{1}, \ldots, \lambda_{d f}\right)$ in terms of matrix product iterations. That analytical representation allows us to check numerically the validity of Inequalities (12) for given $d f, \lambda_{j}, \nu_{j}$, and $m=1, \ldots, M$, for a large bound $M$ (see Appendix E). In Appendix E, when Inequalities (13) are met, we always conclude to the MLR property in the numerical checks as predicted by the theory of Proposition 7. There, we also provide numerical evidence that the domain of validity of the MLR property is relevant for our empirical application. The sufficient conditions (12) and (13) in Proposition 7 yielding the monotone property of density ratios have potentially broad application outside the current setting to show AUMPI properties of other tests based on an asymptotic distribution characterized by a positive definite quadratic form in normal vectors.

## 6 Empirical application

In this section, we test hypotheses about the number of latent factors driving stock returns in short subperiods of the Center for Research in Securities Prices (CRSP) panel. Then, we decompose the cross-sectional variance into systematic and idiosyncratic components. We also check whether there is spanning between the estimated latent factors and standard observed factors.

### 6.1 Testing for the number of latent factors

We consider monthly returns of U.S. common stocks trading on the NYSE, AMEX or NASDAQ between January 1963 and December 2021, and having a non-missing Standard Industrial Classification (SIC) code. We partition subperiods into bull and bear market phases according to the classification methodology of Lunde and Timmermann (2004). ${ }^{19}$ We implement the tests using a rolling window of $T=20$ months, moving forward 12 months each time (adjacent windows over-

[^16]lap by 8 months), thereby ensuring that we can test up to 14 latent factors in each subperiod. The size of the cross-section $n$ ranges from 1768 to 6142 , and the median is 3680 . We only consider stocks with available returns over the whole subperiod, so that our panels are balanced. In each subperiod, we sequentially test $H_{0}(k)$ v.s. $H_{1}(k)$, for $k=0, \ldots, k_{\max }$, where $k_{\max }=14$ is the largest nonnegative integer such that $d f>0$ (see Table 3 in the SMC). We compute the variancecovariance estimator $\hat{\Omega}_{\bar{Z}^{*}}$ using a block structure implied by the partitioning of stocks by the first two digits of their SIC code. The number of blocks ranges from 61 to 87 over the sample, and the number of stocks per block ranges from 1 to 641 . The median number of blocks is 76 and the median number of stocks per block is 21 . We display the p-values of the statistic $L R(k)$ for each subperiod in the upper panel of Figure 1, stopping at the smallest $k$ such that $H_{0}(k)$ is not rejected at level $\alpha_{n}=10 / n_{\max }$, where $n_{\max }$ is the largest cross-sectional sample size over all subperiods, so that $\alpha_{n}=0.16 \%$ in our data. If no such $k$ is found then p -values are displayed up to $k_{\max }$. The $n$-dependent size adjustment controls for the over-rejection problem induced by sequential testing (see Section 6.2 below). Overall, the results point to a higher number of latent factors during bear market phases compared to bull market phases and a decrease of the number of factors over time. ${ }^{20}$ It remains true for the three-month recession periods 1987/09-1987/11 and 2020/01-2020/03, which represent only a fraction of their respective subperiods, although there are "bull" market periods finding a similar number of latent factors. In particular, our results based on a fixed $T$ and large $n$ approach contradict the common wisdom of a single factor model during market downturns due to estimated correlations between equities approaching 1 . It is consistent with the presence of risk factors, such as tail risk or liquidity risk, only showing in stress periods. A rise in the estimated $k$ often happens towards the end of the recession periods. It is consistent with the methodology of Lunde and Timmermann (2004) being early in detecting bear periods (early warning system). The

[^17]results with statistic $\mathscr{T}(k)$ are similar and not reported. The average estimated number of factors is around 7, close to the 4 to 6 factors found by PCA in Bai and Ng (2006) on large time spans of individual stocks. ${ }^{21}$

### 6.2 Decomposing the cross-sectional variance

Building on the results in Pötscher (1983), we can obtain a consistent estimator of the number of latent factors in each subperiod by allowing the asymptotic size $\alpha$ go to zero as $n \rightarrow \infty$ in the sequential testing procedure. We let $\hat{k}$ be defined as the smallest nonnegative integer $k$ satisfying $\operatorname{pval}(k)>\alpha_{n}$, where $\operatorname{pval}(k)$ is the p -value from testing $H_{0}(k)$, and $\alpha_{n}$ is a sequence in $[0,1]$ with $\alpha_{n} \rightarrow 0$. In practice, we take $\alpha_{n}=10 / n_{\max }$. ${ }^{22}$ If no such $k$ is found after sequentially testing $H_{0}(k)$, for $k=0, \ldots, k_{\max }$ at level $\alpha_{n}$, then we take $\hat{k}=k_{\max }+1$. We use the estimate $\hat{k}$ at each subperiod to decompose the path of the cross-sectional variance of stock returns into its systematic and idiosyncratic parts: $\hat{V}_{y, t t}=\hat{F}_{t}^{\prime} \hat{F}_{t}+\hat{V}_{\varepsilon, t t}$, where $\hat{F}$ and $\hat{V}_{\varepsilon}$ are the FA estimates obtained by extracting $\hat{k}$ latent factors. The condition (FA1) ensures that the decomposition holds for any $t$. Such a decomposition is invariant to the choice of normalization for the latent factors. If we look at time averages on a subperiod, we get the decomposition $\overline{\hat{V}}_{y}=\overline{\hat{F}^{\prime} \hat{F}}+\overline{\hat{V}}_{\varepsilon}$, where the overline indicates averaging $\hat{V}_{y, t t}=\hat{F}_{t}^{\prime} \hat{F}_{t}+\hat{V}_{\varepsilon, t t}$ on $t$. In the lower panels of Figure 1, the blue dots correspond to the square root of those quantities for the volatilities, while the ratios $\hat{R}^{2}=\overline{\hat{F}^{\prime}} \hat{F} / \overline{\hat{V}}_{y}$ and $\hat{R}^{2}$ under a single-factor model in the two last panels give measures of goodness-of-fit. ${ }^{23}$ We

[^18]can observe an uptrend in total and idiosyncratic volatilities, while the systematic volatility appears to remain stable over time even if the number of factors has overall decreased over time. ${ }^{24}$ As a result, $\hat{R}^{2}$ is lower on average after the year 2000, indicating a more noisy environment. During the 2007-2008 financial crisis, we can observe a rise in systematic volatility, causing $\hat{R}^{2}$ to reach $59 \%$ during that period. In bear markets, $\hat{R}^{2}$ is often higher. It means that over a bear subperiod, the systematic risk explains a large part of the cross-sectional total variance even if it is not driven by a single factor as reported in Section 6.1. The lowest panel in Figure 1 also signals that $\hat{R}^{2}$ under the constraint of a single-factor model can be way below the one given by the multifactor model. It also means that the idiosyncratic volatility is overestimated if we use a single latent factor only. The plots of the equal-weighted market and firm volatilities used as measures of total and idiosyncratic volatility from a CAPM decomposition in Campbell et al. (2023) show similar patterns as our panels in Figure 1. ${ }^{25}$ Section 4 of Campbell et al. (2023) discusses economic forces (firm fundamentals and investor sentiments) driving the observed time-series variation in average idiosyncratic volatility.

### 6.3 Spanning with observed factors

As discussed in Bai and Ng (2006), we get economic interpretation of latent factors with observed factors when we have spanning between the latent factors and the observed factors to be used as proxies in asset pricing (Shanken (1992)). When $n$ and $T$ are large, Bai and Ng (2006) exploit the asymptotic normality of the empirical canonical correlations between the two sets of factors to investigate spanning under a symmetric role of the two sets. When $T$ is fixed, we suggest the following strategy based on testing for the rank of a matrix. Let us consider $k^{O} \geq k$ empirical

[^19]factors that are excess returns of portfolios ${ }^{26}$, and let $\hat{F}^{O}$ denote the $T \times k^{O}$ matrix of their values with row $t$ given by the transpose of $\hat{f}_{t}^{O}=\frac{1}{n} \sum_{i=1}^{n}\left(y_{i, t}-r_{f, t}\right) z_{i, t}$, where $\frac{1}{n} z_{i, t}$ is a $k^{O} \times 1$ vector of time-varying portfolio weights (long or short positions) based on stocks characteristics. Let matrix $F^{O}$ with rows $f_{t}^{O}=\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{i=1}^{n} E\left[\left(y_{i, t}-r_{f, t}\right) z_{i, t}\right]$ be the corresponding large- $n$ population limit. The notation $\hat{F}^{O}$ makes clear that the sample average of weighted excess returns is an estimate of the population values $F^{O}$. We need to take this into account in the asymptotic analysis of the rank test statistics when $n \rightarrow \infty$. From the factor model under NA, $y_{i, t}=r_{f, t}+f_{t}^{\prime} \tilde{\beta}_{i}+\varepsilon_{i, t}$ (see footnote 2 ), and assuming cross-sectional non-correlation of idiosyncratic errors and portfolio weights, we get $F^{O}=F \Phi^{\prime}$, where $\Phi=\lim _{n \rightarrow \infty} \frac{1}{n} \sum_{i=1}^{n} E\left[z_{i, t}\right] \tilde{\beta}_{i}^{\prime}$ is assumed independent of $t, t=1, \ldots, T$. Hence, the range of $F^{O}$ is a subset of the range of $F$, namely the latent factors span the observed factors (in the population limit sense) by construction. Moreover, $\operatorname{Rank}\left(F^{O}\right) \leq k$. We can test the null hypothesis that $F$ and $F^{O}$ span the same linear spaces, namely matrices $F$ and $F^{0}$ have the same range. Such a null hypothesis is equivalent to the rank condition: $\operatorname{Rank}\left(F^{O}\right)=k$.

We build on the rank testing literature; see e.g. Cragg and Donald (1996), Robin and Smith (RS, 2000), Kleibergen and Paap (KP, 2006), Al-Sadoon (2017). ${ }^{27}$ We use in particular the RS and KP statistics. For those tests, the null hypothesis is that a given matrix has a reduced rank $r$ against the alternative that the rank is greater than $r$. Hence, to test for spanning by the empirical factors, we consider the null hypothesis $H_{0, s p}(r): \operatorname{Rank}\left(F^{O}\right)=r$ against the alternative $H_{1, s p}(r)$ : $\operatorname{Rank}\left(F^{O}\right)>r$, for any integer $r<k .{ }^{28}$ We use the asymptotic expansion $\hat{F}^{O}=\tilde{F}^{O}+\frac{1}{\sqrt{n}} \Psi_{F^{O}, n}$, where $\tilde{F}^{O}=F \Phi_{n}^{\prime}$ with $\Phi_{n}=\frac{1}{n} \sum_{i=1}^{n} E\left[z_{i, t}\right] \tilde{\beta}_{i}^{\prime}$, and the rows of matrix $\Psi_{F^{0}, n}$ are given by $\Psi_{F, n, t}=$ $\frac{1}{\sqrt{n}} \sum_{i=1}^{n}\left(\eta_{i, t} f_{t}+\varepsilon_{i, t} z_{i, t}\right)$ with $\eta_{i, t}:=\left(z_{i, t}-E\left[z_{i, t}\right]\right) \tilde{\beta}_{i}^{\prime}$. Under $H_{0, s p}(r)$, we assume that $\Phi_{n}$ has

[^20]the same null space as $\Phi$, in particular $\Phi_{n}$ has rank $r$, for $n$ large enough. ${ }^{29}$ We assume the CLT $\operatorname{vec}\left(\Psi_{F^{0}, n}\right) \Rightarrow N\left(0, \Omega_{\Psi}\right)$. Further, we use the Singular Value Decomposition (SVD) of matrix $\hat{F}^{O}=\hat{U} \hat{\mathcal{S}} \hat{V}^{\prime}$. Then, the RS and KP statistics are the quadratic forms $\mathscr{S}_{R S}=\operatorname{nvec}\left(\hat{\mathcal{S}}_{22}\right)^{\prime} \operatorname{vec}\left(\hat{\mathcal{S}}_{22}\right)$ and $\mathscr{S}_{K P}=\operatorname{nvec}\left(\hat{\mathcal{S}}_{22}\right)^{\prime} \hat{\Omega}_{\mathcal{S}}^{-1} \operatorname{vec}\left(\hat{\mathcal{S}}_{22}\right)$, where $\hat{\mathcal{S}}_{22}$ is the lower-right $(T-r) \times\left(k^{O}-r\right)$ block of matrix $\hat{\mathcal{S}}$. Here, $\hat{\Omega}_{\mathcal{S}}=\left(\hat{V}_{k{ }^{O_{-r}}} \otimes \hat{U}_{T-r}\right)^{\prime} \hat{\Omega}_{\Psi}\left(\hat{V}_{k o_{-r}} \otimes \hat{U}_{T-r}\right)$, where $\hat{U}_{T-r}$ and $\hat{V}_{k 0_{-r}}$ are the $T \times(T-r)$ and $k^{O} \times\left(k^{O}-r\right)$ matrices in the block forms $\hat{U}=\left[\hat{U}_{r}: \hat{U}_{T-r}\right]$ and $\hat{V}=\left[\hat{V}_{r}: \hat{V}_{k} O_{-}\right]$. In the SMC, we design a consistent estimator $\hat{\Omega}_{\Psi}$ of $\Omega_{\Psi}$ building on a block structure for the characteristics akin to Assumption 2 and a stationarity condition. The definitions of the test statistics $\mathscr{S}_{R S}$ and $\mathscr{S}_{K P}$ are equivalent to those in the original RS and KP papers. The asymptotic distributions under $H_{0, s p}(r)$ are $\mathscr{S}_{R S} \Rightarrow \sum_{j=1}^{(T-r)\left(k^{O}-r\right)} \delta_{j}\left(\Omega_{\mathcal{S}}\right) \chi_{j}^{2}(1)$ and $\mathscr{S}_{K P} \Rightarrow \chi^{2}\left[(T-r)\left(k^{O}-r\right)\right]$, where $\Omega_{\mathcal{S}}=$ $\left(V_{k} O_{-r} \otimes U_{T-r}\right)^{\prime} \Omega_{\Psi}\left(V_{k O_{-r}} \otimes U_{T-r}\right)$ is assumed non-singular.

We build the empirical matrix $\hat{F}^{O}$ with the time-varying portfolio weights of the Fama-French five-factor model (Fama and French (2015)) plus the momentum factor (Carhart (1997)), i.e., $k^{O}=$ 6. In the two panels of Figure 2, we can observe that the rank tests point most of the time at a low reduced rank $r$ either 1 or 2 , with only occasionally 3 or 4 , for the matrix $\hat{F}^{O}$. Observed factors struggle spanning latent factors since their associated linear space is of a dimension smaller than the one of the latent factor space. The discrepancy between the dimensions of the two factor spaces has decreased over time. According to the KP statistic, the rank deficiency of $\hat{F}^{O}$ is often less pronounced in bear markets indicating less redundancy between the observed factors.

## 7 Concluding remarks

In this paper, we develop a new theory of Factor Analysis in short panels beyond the Gaussian and i.i.d. cases. We establish the AUMPI property of the LR statistic for testing hypotheses on the number of latent factors. Our results for short subperiods of the CRSP panel of US stock returns contradict the common wisdom of a single factor during market downturns. In bear markets,

[^21]systematic risk explains a large part of the cross-sectional variance, and is not spanned by traditional empirical factors. Our new methodology can be used to address relevant empirical questions in applications beyond asset pricing. For example, in analysis of education when a panel consists of students repeatedly tested along different cognitive domains in mathematics and science (Freyberger (2018)) or interviewed in successive waves (Sarzosa and Urzúa (2021)), in analysis of particular time spans in long panel data of wages (Gobillon, Magnac, and Roux (2022)) or in analysis of unemployment for a panel of counties followed on a couple of years (Hagedorn, Manovskii, and Mitman (2015)).

## References

Ahn, S., and Horenstein, A., 2013. Eigenvalue ratio test for the number of factors. Econometrica 81 (3), 1203-1227.

Ahn, S., Horenstein, A., and Wang, N., 2018. Beta matrix and common factors in stock returns. Journal of Financial and Quantitative Analysis 53 (3), 1417-1440.

Aigner, D., Hsiao, C., Kapteyn A., and Wansbeek, T., 1984. Latent variable models in econometrics, in Handbook of Econometrics, Volume II, Z. Griliches and M.D. Intriligator Eds., 1321-1393. Aït-Sahalia, Y., and Jacod, J., 2014. High-frequency financial econometrics. Princeton University Press.

Al-Sadoon, M., 2017. A unifying theory of tests of rank. Journal of Econometrics 199 (1), 49-62.
Andersen, T., Bollerslev, T., Diebold, F., and Labys, P., 2003. Modeling and forecasting realized volatility. Econometrica 71(2), 579-625.

Anderson, T. W., 2003. An introduction to multivariate statistical analysis. Wiley.
Anderson, T. W., and Rubin, H., 1956. Statistical inference in factor analysis. Proceedings of the Third Berkeley Symposium in Mathematical Statistics and Probability 5, 11-150.

Anderson, T. W. and Amemiya, Y., 1988. The asymptotic normal distribution of estimators in factor analysis under general conditions. Annals of Statistics 16 (2), 759-771.

Ando, T., and Bai, J., 2015. Asset pricing with a general multifactor structure. Journal of Financial Econometrics 13 (3), 556-604.

Bai, J., 2003. Inferential theory for factor models of large dimensions. Econometrica 71 (1), 135-171.

Bai, J., 2009. Panel data models with interactive effects. Econometrica, 77 (4), 1229-1279.
Bai, J., and Li, K., 2012. Statistical analysis of factor models of high dimension. Annals of Statistics 40 (1), 436-465.

Bai, J., and Li, K., 2016. Maximum likelihood estimation and inference for approximate factor models of high dimension. Review of Economics and Statistics 98 (2), 298-309.

Bai, J., and Ng, S., 2002. Determining the number of factors in approximate factor models. Econometrica 70 (1), 191-221.

Bai, J., and Ng, S., 2006. Evaluating latent and observed factors in macroeconomics and finance. Journal of Econometrics 131 (1), 507-537.

Barigozzi, M., and Hallin, M., 2016. Generalized dynamic factor models and volatilities: recovering the market volatility shocks. Econometrics Journal 19 (1), 33-60.

Barndorff-Nielsen, O., and Shephard, N., 2002. Econometric analysis of realized volatility and its use in estimating stochastic volatility models. Journal of the Royal Statistical Society, Series B, 64 (2), 253-280.

Bell, E. T., 1934. Exponential polynomials. Annals of Mathematics 35 (2), 258-277.
Bernstein, D., 2009. Matrix mathematics: Theory, facts and formulas. Princeton University Press.
Bickel, P. J., and Levina, E., 2008. Covariance regularization by thresholding. Annals of Statistics 36 (6), 2577-2604.

Campbell, J., Lettau, M., Malkiel, B., and Xu, Y., 2023. Idiosyncratic equity risk two decades later. Critical Finance Review, 12, 203-223.

Caner, M., and Han, X., 2014. Selecting the correct number of factors in approximate factor models: the large panel case with group bridge estimator. Journal of Business and Economic Statistics 32 (3), 359-374.

Carhart, M., 1997. On persistence in mutual fund performance. Journal of Finance. 52 (1), 57-82. Chamberlain, G., 1992. Efficiency bounds for semi-parametric regression. Econometrica 60 (3), 567-596.

Cochrane, J., 2005. Asset pricing. Princeton University Press.
Connor, G., and Korajczyk, R., 1986. Performance measurement with the arbitrage pricing theory: A new framework for analysis. Journal of Financial Economics 15 (3), 373-394.

Connor, G., and Korajczyk, R., 1993. A test for the number of factors in an approximate factor model. Journal of Finance 48 (4), 1263-1291.

Cragg, J., and Donald, S., 1996. On the asymptotic properties of LDU-based tests of the rank of a matrix. Journal of the American Statistical Association 91 (435), 1301-1309.

Eaton, M., 1987. Multivariate statistics. A vector space approach. Institute of Mathematical Statistics Lecture notes-Monograph series, Vol. 53.

Engle, R., 1984. Wald, likelihood ratio, and Lagrange multiplier tests in econometrics, in Handbook of Econometrics, Volume II, Z. Griliches and M.D. Intriligator Eds., 775-826.

Fama, E., and French, K., 2015. A five-factor asset pricing model. Journal of Financial Economics 116 (1), 1-22.

Fortin, A.-P., Gagliardini, P., and Scaillet, O., 2022. Eigenvalue tests for the number of latent factors in short panels. Journal of Financial Econometrics, forthcoming.
Freyberger, J., 2018. Non-parametric panel data models with interactive fixed effects. Review of Economic Studies 85 (3), 1824-1851.

Gagliardini, P., Ossola, E., and Scaillet, O., 2016. Time-varying risk premium in large crosssectional equity datasets. Econometrica 84 (3), 985-1046.

Gagliardini, P., Ossola, E., and Scaillet, O., 2019. A diagnostic criterion for approximate factor structure. Journal of Econometrics 21 (2), 503-521.

Gagliardini, P., Ossola, E., and Scaillet, O., 2020. Estimation of large dimensional conditional factor models in finance, in Handbook of Econometrics, Volume 7A, S. Durlauf, L. Hansen, J. Heckman, and R. Matzkin Eds., 219-282.

Gobillon, L., Magnac, T., and Roux, S., 2022. Lifecycle wages and human capital investments: Selection and missing data. TSE working paper.

Hagedorn, M., Manovskii, I., and Mitman, K., 2015. The impact of unemployment benefit extensions on employment: The 2014 employment miracle. NBER working paper.

Kapetanios, G., 2010. A testing procedure for determining the number of factors in approximate factor models with large datasets. Journal of Business and Economic Statistics 28 (3), 397-409.

Kim, S., and Skoulakis, G., 2018. Ex-post risk premia estimation and asset pricing tests using large cross-sections: the regression-calibration approach. Journal of Econometrics 204 (2), 159-188. Kleibergen, F., and Paap, R., 2006. Generalized reduced rank tests using the singular value decomposition. Journal of Econometrics 133 (1), 97-126.

Kleibergen, F., and Zhan, Z., 2023. Identification-robust inference for risk premia in short panels. Working paper, University of Amsterdam.

Kotz, S., Johnson, N., and Boyd, D., 1967. Series representations of distributions of quadratic forms in Normal variables II. Non-central case. Annals of Mathematical Statistics 38 (3), 838-848. Lehmann, E., and Romano, D., 2005. Testing statistical hypotheses, Springer Texts in Statistics. Lunde, A., and Timmermann, A., 2004. Duration dependence in stock prices: An analysis of bull and bear markets. Journal of Business and Economic Statistics 22 (3), 253-273.

Magnus, J., and Neudecker, H., 2007. Matrix differential calculus, with applications in statistics and econometrics. Wiley.

Miravete, E., 2011. Convolution and composition of totally positive random variables in economics. Journal of Mathematical Economics 47 (4), 479-490.

Onatski, A., 2009. Testing hypotheses about the number of factors in large factor models. Econometrica 77 (5), 1447-1479.

Onatski, A., 2010. Determining the number of factors from empirical distribution of eigenvalues. Review of Economics and Statistics 92 (4), 1004-1016.

Onatski, A. 2022. Comment on "Eigenvalue tests for the number of latent factors in short panels" by A.-P. Fortin, P. Gagliardini and O. Scaillet, Journal of Financial Econometrics, forthcoming.

Pötscher, B., 1983. Order estimation in ARMA-models by Lagrangian multiplier tests, Annals of Statistics 11 (3), 872-885.

Raponi, V., Robotti, C., and Zaffaroni, P., 2020. Testing beta-pricing models using large crosssections. Review of Financial Studies 33 (6), 2796-2842.

Renault, E., Van Der Heijden, T., and Werker, B., 2022. Arbitrage pricing theory for idiosyncratic variance factors. Journal of Financial Econometrics, forthcoming.

Robin, J.-M., and Smith, R., 2000. Tests of rank. Econometric Theory 16 (2), 151-175.
Sarzosa, M., and Urzúa, S., 2021. Bullying among adolescents: The role of skills. Quantitative Economics 12 (3), 945-980.

Shanken, J. 1992. On the estimation of beta pricing models. Review of Financial Studies 5 (1), 1-33.

Stock, J., and Watson, M., 2002. Forecasting using principal components from a large number of predictors. Journal of the American Statistical Association 97 (460), 1167-1179.

Tao, T. 2012. Topics in random matrix theory. Graduate Studies in Mathematics, Volume 132, American Mathematical Society.

White, H., 1982. Maximum likelihood estimation of misspecified models. Econometrica 50 (1), 1-25.

Zaffaroni, P., 2019. Factor models for conditional asset pricing. Working Paper.

## Appendix

## A Regularity assumptions

In this appendix, we list and comment the additional assumptions used to derive the large sample properties of the estimators and test statistics. We often denote by $C>0$ a generic constant. Set $\Theta$ is a compact subset of $\left\{\theta=\left(\operatorname{vec}(F)^{\prime}, \operatorname{diag}\left(V_{\varepsilon}\right)^{\prime}\right)^{\prime} \in \mathbb{R}^{r} \quad: \quad V_{\varepsilon}\right.$ is diagonal and positive definite, $F^{\prime} V_{\varepsilon}^{-1} F$ is diagonal, with diagonal elements ranked in decreasing order $\}$ with $r=(T+$

1) $k$, and function $L_{0}(\theta)=-\frac{1}{2} \log |\Sigma(\theta)|-\frac{1}{2} \operatorname{Tr}\left(V_{y} \Sigma(\theta)^{-1}\right)$ is the population FA criterium, where $\Sigma(\theta)=F F^{\prime}+V_{\varepsilon}$ and $V_{y}=\operatorname{plim}_{n \rightarrow \infty} \hat{V}_{y}$. Further, $\theta_{0}=\left(\operatorname{vec}\left(F_{0}\right)^{\prime}, \operatorname{diag}\left(V_{\varepsilon}^{0}\right)^{\prime}\right)^{\prime}$ denotes the vector of true parameter values under $H_{0}(k)$ and is an interior point of set $\Theta$.

Assumption A. 1 The loadings are normalized such that $\bar{\beta}=\frac{1}{n} \sum_{i=1}^{n} \beta_{i}=0$ and $\tilde{V}_{\beta}:=\frac{1}{n} \sum_{i=1}^{n} \beta_{i} \beta_{i}^{\prime}$ $=I_{k}$, for any $n$. Moreover, $\left|\beta_{i}\right| \leq C$, for all $i$.

Assumption A. 2 We have $E\left[w_{i, t}^{8}\right] \leq C$ and $\left|\sigma_{i, j}\right| \leq C$, for all $i, j, t$.
Assumption A. 3 Under the null hypothesis $H_{0}(k)$, we have: $\Sigma(\theta)=\Sigma\left(\theta_{0}\right), \theta \in \Theta \Rightarrow \theta=\theta_{0}$, up to sign changes in the columns of $F$.

Assumption A. 4 (a) The $\frac{T(T+1)}{2} \times \frac{T(T+1)}{2}$ symmetric matrix $D=\lim _{n \rightarrow \infty} D_{n}$ exists, where $D_{n}=$ $\frac{1}{n} \sum_{i=1}^{n} \sigma_{i i}^{2} V\left[\operatorname{vech}\left(w_{i} w_{i}^{\prime}\right)\right]$. (b) We have $\delta_{T(T+1) / 2}\left(V\left[\operatorname{vech}\left(w_{i} w_{i}^{\prime}\right)\right]\right) \geq \underline{c}$, for all $i \in \bar{S}$, where $\bar{S} \subset\{1, \ldots, n\}$ with $\frac{1}{n} \sum_{i=1}^{n} 1_{i \in \bar{S}} \geq 1-\frac{1}{2 \bar{C}}$, for constants $\bar{C}, \bar{c}>0$, such that $\sigma_{i i} \leq \bar{C}$. (c) We have $\lim _{n \rightarrow \infty} \kappa_{n}=\kappa$ for a constant $\kappa \geq 0$, where $\kappa_{n}:=\frac{1}{n} \sum_{m=1}^{J_{n}}\left(\sum_{i \neq j \in I_{m}} \sigma_{i j}^{2}\right)$.
Assumption A. 5 Under the alternative hypothesis $H_{1}(k)$, (a) function $L_{0}(\theta)$ has a unique maximizer $\theta^{*}=\left(\operatorname{vec}\left(F^{*}\right)^{\prime} \text {, } \operatorname{diag}\left(V_{\varepsilon}^{*}\right)^{\prime}\right)^{\prime}$ over $\Theta$, and $(b)$ we have $V_{y} \neq F F^{\prime}+V_{\varepsilon}$, for any $T \times k$ matrix $F$ and any diagonal positive definite matrix $V_{\varepsilon}$.
Assumption A.6 Matrix $Q_{\beta}:=\lim _{n \rightarrow \infty} \frac{1}{n} \beta^{\prime} \Sigma \beta$ is positive definite.
Assumptions A. 1 and A. 2 require uniform bounds on factor loadings as well as on covariances and higher-order moments of the idiosyncratic errors. Assumption A. 3 implies global identification in the FA model (see Lemma 5). Assumptions A.1-A. 3 yield consistency of FA estimators (see proof of Lemma 6). We use Assumption A. 4 together with Assumption A. 2 to invoke a CLT based on a multivariate Lyapunov condition (see proof of Lemma 1) to establish the asymptotic distribution of the test statistics. To ease the verification of the Lyapunov condition, we bound a fourth-order moment of squared errors, which explains why we require finite eight-order moments in Assumption A.2. We could relax this condition at the expense of a more sophisticated proof of Lemma 1. The mild Assumption A. 4 (b) requires that the smallest eigenvalue of $V\left[\operatorname{vech}\left(w_{i} w_{i}^{\prime}\right)\right]$ is bounded away from 0 for all assets $i$ up to a small fraction. In Assumption A. 4 (c), in order
to have $\kappa_{n}$ bounded, we need either mixing dependence in idiosyncratic errors within blocks, i.e., $\left|\sigma_{i, j}\right| \leq C \rho^{|i-j|}$ for $i, j \in I_{m}$ and $0 \leq \rho<1$, or vanishing correlations, i.e., $\left|\sigma_{i, j}\right| \leq C b_{m, n}^{-\bar{s}}$ for all $i \neq j \in I_{m}$ and a constant $\bar{s} \geq 1 / 2$, with blocks of equal size. In Assumption A.5, part (a) defines the pseudo-true parameter value (White (1982)) under the alternative hypothesis, and part (b) is used to establish the consistency of the LR test under global alternatives (see proof of Proposition 4). Finally, Assumption A. 6 is used to apply a Lyapunov CLT (see proof of Lemma 8) when deriving the asymptotic normality of the FA estimators.

## B Proofs of Propositions 1-7

Proof of Proposition 1: Let $\hat{U}$ be the $T \times k$ matrix whose orthonormal columns are the eigenvectors for the $k$ largest eigenvalues of matrix $\hat{V}_{\varepsilon}^{-1 / 2} \hat{V}_{y} \hat{V}_{\varepsilon}^{-1 / 2}$. Those eigenvalues are $1+\hat{\gamma}_{j}$, $j=1, \ldots, k$, while it holds $\hat{F}=\hat{V}_{\varepsilon}^{1 / 2} \hat{U} \hat{\Gamma}^{1 / 2}$, where $\hat{\Gamma}=\operatorname{diag}\left(\hat{\gamma}_{1}, \ldots, \hat{\gamma}_{k}\right)$. We have $I_{T}-\hat{U} \hat{U}^{\prime}=I_{T}-$ $\hat{V}_{\varepsilon}^{-1 / 2} \hat{F} \hat{\Gamma}^{-1} \hat{F}^{\prime} \hat{V}_{\varepsilon}^{-1 / 2}=I_{T}-\hat{V}_{\varepsilon}^{-1 / 2} \hat{F}\left(\hat{F}^{\prime} \hat{V}_{\varepsilon}^{-1} \hat{F}\right)^{-1} \hat{F}^{\prime} \hat{V}_{\varepsilon}^{-1 / 2}=\hat{V}_{\varepsilon}^{-1 / 2} M_{\hat{F}, \hat{V}_{\varepsilon}} \hat{V}_{\varepsilon}^{1 / 2}=\hat{V}_{\varepsilon}^{1 / 2} M_{\hat{F}, \hat{V}_{\varepsilon}}^{\prime} \hat{V}_{\varepsilon}^{-1 / 2}$. Thus, $\hat{S}=\left(I_{T}-\hat{U} \hat{U}^{\prime}\right)\left(\hat{V}_{\varepsilon}^{-1 / 2} \hat{V}_{y} \hat{V}_{\varepsilon}^{-1 / 2}-I_{T}\right)\left(I_{T}-\hat{U} \hat{U}^{\prime}\right)$. By the spectral decomposition of $\hat{V}_{\varepsilon}^{-1 / 2} \hat{V}_{y} \hat{V}_{\varepsilon}^{-1 / 2}$, we get $\left(I_{T}-\hat{U} \hat{U}^{\prime}\right)\left(\hat{V}_{\varepsilon}^{-1 / 2} \hat{V}_{y} \hat{V}_{\varepsilon}^{-1 / 2}-I_{T}\right)\left(I_{T}-\hat{U} \hat{U}^{\prime}\right)=\sum_{j=k+1}^{T} \hat{\gamma}_{j} \hat{P}_{j}$, where the $\hat{P}_{j}$ are the orthogonal projection matrices onto the eigenspaces for the $T-k$ smallest eigenvalues. Then, Part (a) follows. Part (b) is a consequence of the squared Frobenius norm of a symmetric matrix being equal to the sum of its squared eigenvalues. For Part (c), let $P_{\hat{F}, \hat{V}_{\varepsilon}}=$ $I_{T}-M_{\hat{F}, \hat{V}_{\varepsilon}}$ and note that $\hat{F} \hat{F}^{\prime}=P_{\hat{F}, \hat{V}_{\varepsilon}}\left(\hat{V}_{y}-\hat{V}_{\varepsilon}\right)+\left(\hat{V}_{y}-\hat{V}_{\varepsilon}\right) P_{\hat{F}, \hat{V}_{\varepsilon}}^{\prime}-P_{\hat{F}, \hat{V}_{\varepsilon}}\left(\hat{V}_{y}-\hat{V}_{\varepsilon}\right) P_{\hat{F}, \hat{V}_{\varepsilon}}^{\prime}$ $=\hat{V}_{y}-\hat{V}_{\varepsilon}-M_{\hat{F}, \hat{V}_{\varepsilon}}\left(\hat{V}_{y}-\hat{V}_{\varepsilon}\right) M_{\hat{F}, \hat{V}_{\varepsilon}}^{\prime}$, where the first equality is because the three terms on the RHS are all equal to $\hat{F} \hat{F}^{\prime}$ by (FA2). The conclusion follows from (FA1) and $\hat{V}_{\varepsilon}$ being diagonal. Finally, Part (d) follows because $\frac{1}{n} \hat{\varepsilon} \hat{\varepsilon}^{\prime}=M_{\hat{F}, \hat{V}_{\varepsilon}} \hat{V}_{y} M_{\hat{F}, \hat{V}_{\varepsilon}}^{\prime}, \hat{V}_{\varepsilon} M_{\hat{F}, \hat{V}_{\varepsilon}}^{\prime}=M_{\hat{F}, \hat{V}_{\varepsilon}} \hat{V}_{\varepsilon}$, and $M_{\hat{F}, \hat{V}_{\varepsilon}}$ is idempotent, which implies $\hat{V}_{\varepsilon}^{-1 / 2} M_{\hat{F}, \hat{V}_{\varepsilon}} \hat{V}_{\varepsilon} M_{\hat{F}, \hat{V}_{\varepsilon}}^{\prime} \hat{V}_{\varepsilon}^{-1 / 2}=\hat{V}_{\varepsilon}^{-1 / 2} M_{\hat{F}, \hat{V}_{\varepsilon}} \hat{V}_{\varepsilon}^{1 / 2}$.

Proof of Proposition 2: Let us substitute $\hat{V}_{y}=F F^{\prime}+\tilde{V}_{\varepsilon}+\frac{1}{\sqrt{n}} \Psi_{y}+o_{p}\left(\frac{1}{\sqrt{n}}\right)$ into (FA2) and rearrange to obtain $\hat{F} \hat{\Gamma}-F F^{\prime} \hat{V}_{\varepsilon}^{-1} \hat{F}=\frac{1}{\sqrt{n}} \Psi_{y} \hat{V}_{\varepsilon}^{-1} \hat{F}+\left(\tilde{V}_{\varepsilon} \hat{V}_{\varepsilon}^{-1}-I_{T}\right) \hat{F}+o_{p}\left(\frac{1}{\sqrt{n}}\right)$. From $\hat{V}_{\varepsilon}=\tilde{V}_{\varepsilon}+\frac{1}{\sqrt{n}} \Psi_{\varepsilon}+o_{p}\left(\frac{1}{\sqrt{n}}\right)$, we have $\tilde{V}_{\varepsilon} \hat{V}_{\varepsilon}^{-1}-I_{T}=-\frac{1}{\sqrt{n}} \Psi_{\varepsilon} \hat{V}_{\varepsilon}^{-1}+o_{p}\left(\frac{1}{\sqrt{n}}\right)$. Substituting into
the above equation and right multiplying both sides by $\left(F^{\prime} \hat{V}_{\varepsilon}^{-1} \hat{F}\right)^{-1}$ gives $\hat{F} \hat{\mathcal{D}}-F=\frac{1}{\sqrt{n}}\left(\Psi_{y}-\right.$ $\left.\Psi_{\varepsilon}\right) \hat{V}_{\varepsilon}^{-1} \hat{F}\left(F^{\prime} \hat{V}_{\varepsilon}^{-1} \hat{F}\right)^{-1}+o_{p}\left(\frac{1}{\sqrt{n}}\right)$, where $\hat{\mathcal{D}}:=\hat{\Gamma}\left(F^{\prime} \hat{V}_{\varepsilon}^{-1} \hat{F}\right)^{-1}$. By the root- $n$ convergence of the FA estimates (see Section D.4.1), we get

$$
\begin{equation*}
\hat{F} \hat{\mathcal{D}}-F=\frac{1}{\sqrt{n}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F \Gamma^{-1}+o_{p}\left(\frac{1}{\sqrt{n}}\right), \tag{B.1}
\end{equation*}
$$

and $\hat{\mathcal{D}}=I_{k}+O_{p}\left(\frac{1}{\sqrt{n}}\right)$, where $\Gamma=\operatorname{diag}\left(\gamma_{1}, \ldots, \gamma_{k}\right)$. We can push the expansion by plugging into (B.1) the expansion of $\hat{\mathcal{D}}$. We have $F^{\prime} \hat{V}_{\varepsilon}^{-1} \hat{F}=\left[I_{k}-(\hat{F}-F)^{\prime} \hat{V}_{\varepsilon}^{-1} \hat{F} \hat{\Gamma}^{-1}\right] \hat{\Gamma}$, so that $\hat{\mathcal{D}}=$ $\left[I_{k}-(\hat{F}-F)^{\prime} \hat{V}_{\varepsilon}^{-1} \hat{F} \hat{\Gamma}^{-1}\right]^{-1}=I_{k}+(\hat{F}-F)^{\prime} V_{\varepsilon}^{-1} F \Gamma^{-1}+o_{p}\left(\frac{1}{\sqrt{n}}\right)$. By plugging into (B.1), we get:

$$
\begin{equation*}
\hat{F}-F+F\left[(\hat{F}-F)^{\prime} V_{\varepsilon}^{-1} F \Gamma^{-1}\right]=\frac{1}{\sqrt{n}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F \Gamma^{-1}+o_{p}\left(\frac{1}{\sqrt{n}}\right) . \tag{B.2}
\end{equation*}
$$

By multiplying both sides with $M_{F, V_{\varepsilon}}$, we get $M_{F, V_{\varepsilon}}(\hat{F}-F)=\frac{1}{\sqrt{n}} M_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F \Gamma^{-1}+$ $o_{p}\left(\frac{1}{\sqrt{n}}\right)$. Then, $\hat{F}-F=\frac{1}{\sqrt{n}} M_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F \Gamma^{-1}+\frac{1}{\sqrt{n}} F A+o_{p}\left(\frac{1}{\sqrt{n}}\right)$, where $A$ is a random $k \times k$ matrix to be determined next. By plugging into (B.2), we get $F\left(A+A^{\prime}\right)=P_{F, V_{\varepsilon}}\left(\Psi_{y}-\right.$ $\left.\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F \Gamma^{-1}+o_{p}(1)$. By multiplying both sides by $\frac{1}{2} \Gamma^{-1} F^{\prime} V_{\varepsilon}^{-1}$ and using $F^{\prime} V_{\varepsilon}^{-1} P_{F, V_{\varepsilon}}=F^{\prime} V_{\varepsilon}^{-1}$, we get the symmetric part of matrix $A$, i.e., $\frac{1}{2}\left(A+A^{\prime}\right)=\frac{1}{2} \Gamma^{-1} F^{\prime} V_{\varepsilon}^{-1}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F \Gamma^{-1}$ (we include higher-order terms in the remainder $o_{p}\left(\frac{1}{\sqrt{n}}\right)$ ). Thus, $\hat{F}-F=\frac{1}{\sqrt{n}} \Psi_{F}+o_{p}\left(\frac{1}{\sqrt{n}}\right)$, where

$$
\begin{equation*}
\Psi_{F}=M_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F \Gamma^{-1}+\frac{1}{2} P_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F \Gamma^{-1}+F \tilde{A}, \tag{B.3}
\end{equation*}
$$

and $\tilde{A}=\frac{1}{2}\left(A-A^{\prime}\right)$ is an antisymmetric $k \times k$ random matrix. To find the antisymmetric matrix $\tilde{A}=\left(\tilde{a}_{\ell, j}\right)$, we use that $\hat{F}^{\prime} \hat{V}_{\varepsilon}^{-1} \hat{F}$ is diagonal. Plugging the expansions of the FA estimates, for the term at order $1 / \sqrt{n}$ we get that the out-of-diagonal elements of matrix $\Psi_{F}^{\prime} V_{\varepsilon}^{-1} F+F^{\prime} V_{\varepsilon}^{-1} \Psi_{F}-$ $F^{\prime} V_{\varepsilon}^{-1} \Psi_{\varepsilon} V_{\varepsilon}^{-1} F=\frac{1}{2} \Gamma^{-1} F^{\prime} V_{\varepsilon}^{-1}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F+\frac{1}{2} F^{\prime} V_{\varepsilon}^{-1}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F \Gamma^{-1}+\Gamma \tilde{A}-\tilde{A} \Gamma-$ $F^{\prime} V_{\varepsilon}^{-1} \Psi_{\varepsilon} V_{\varepsilon}^{-1} F$ are nil. Setting the $(\ell, j)$ element of this matrix equal to 0 , we get $\tilde{a}_{\ell, j}=-\tilde{a}_{j, \ell}=$ $\frac{1}{\gamma_{j}-\gamma_{\ell}}\left[\frac{1}{2}\left(\frac{1}{\gamma_{j}}+\frac{1}{\gamma_{\ell}}\right) F_{\ell}^{\prime} V_{\varepsilon}^{-1}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F_{j}-F_{\ell}^{\prime} V_{\varepsilon}^{-1} \Psi_{\varepsilon} V_{\varepsilon}^{-1} F_{j}\right]$, for $j \neq \ell$. Then, from Equation (B.3), the $j$ th column of $\Psi_{F}$ is $\Psi_{F_{j}}=\frac{1}{\gamma_{j}} M_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F_{j}+\frac{1}{2 \gamma_{j}} P_{F_{j}, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F_{j}+$ $\sum_{\ell=1: \ell \neq j}^{k} \frac{1}{\gamma_{j}-\gamma_{\ell}} P_{F_{\ell}, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) V_{\varepsilon}^{-1} F_{j}-\sum_{\ell=1: \ell \neq j}^{k} \frac{\gamma_{\ell}}{\gamma_{j}-\gamma_{\ell}} P_{F_{\ell}, V_{\varepsilon}} \Psi_{\varepsilon} V_{\varepsilon}^{-1} F_{j}$, where we use $P_{F, V_{\varepsilon}}=$ $\sum_{\ell=1}^{k} P_{F_{\ell}, V_{\varepsilon}}$. Part (a) follows.

Let us now prove part (b). The asymptotic expansion of condition (FA1) yields:

$$
\begin{equation*}
\operatorname{diag}\left(\Psi_{y}\right)=\operatorname{diag}\left(\sum_{j=1}^{k}\left(F_{j} \Psi_{F_{j}}^{\prime}+\Psi_{F_{j}} F_{j}^{\prime}\right)+\Psi_{\varepsilon}\right) \tag{B.4}
\end{equation*}
$$

From part (a) and the definition of $P_{F_{j}, V_{\varepsilon}}$ we have $\sum_{j=1}^{k} \Psi_{F_{j}} F_{j}^{\prime}=\frac{1}{2} \sum_{j=1}^{k} P_{F_{j}, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) P_{F_{j}, V_{\varepsilon}}^{\prime}+$ $M_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) P_{F, V_{\varepsilon}}^{\prime}+\sum_{\ell \neq j} \frac{\gamma_{j}}{\gamma_{j}-\gamma_{\ell}} P_{F_{\ell}, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) P_{F_{j}, V_{\varepsilon}}^{\prime}-\sum_{\ell \neq j}^{k} \frac{\gamma_{\ell} \gamma_{j}}{\gamma_{j}-\gamma_{\ell}} P_{F_{\ell}, V_{\varepsilon}} \Psi_{\varepsilon} P_{F_{j}, V_{\varepsilon}}^{\prime}=: N_{1}+$ $N_{2}+N_{3}+N_{4}$, where $P_{F, V_{\varepsilon}}=\sum_{j=1}^{k} P_{F_{j}, V_{\varepsilon}}=I_{T}-M_{F, V_{\varepsilon}}$ and $\sum_{\ell \neq j}$ denotes the double sum over $j, \ell=1, \ldots, k$ such that $\ell \neq j$. Matrix $N_{1}$ is symmetric and it contributes $2 N_{1}$ to the RHS of (B.4). Instead, matrix $N_{4}$ is antisymmetric (it can be seen by interchanging indices $j$ and $\ell$ in the summation) and it does not contribute to the RHS of (B.4). For matrix $N_{3}$ we have $N_{3}+$ $N_{3}^{\prime}=\sum_{\ell \neq j} \frac{\gamma_{j}}{\gamma_{j}-\gamma_{\ell}} P_{F_{\ell}, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) P_{F_{j}, V_{\varepsilon}}^{\prime}+\sum_{\ell \neq j} \frac{\gamma_{\ell}}{\gamma_{\ell}-\gamma_{j}} P_{F_{\ell}, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) P_{F_{j}, V_{\varepsilon}}^{\prime}=\sum_{\ell \neq j} P_{F_{\ell}, V_{\varepsilon}}\left(\Psi_{y}-\right.$ $\left.\Psi_{\varepsilon}\right) P_{F_{j}, V_{\varepsilon}}^{\prime}=\sum_{\ell, j} P_{F_{\ell}, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) P_{F_{j}, V_{\varepsilon}}^{\prime}-\sum_{j} P_{F_{j}, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) P_{F_{j}, V_{\varepsilon}}^{\prime}=P_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) P_{F, V_{\varepsilon}}^{\prime}-2 N_{1}$, where we have interchanged $\ell$ and $j$ in the first equality when writing $N_{3}^{\prime}$. Thus, we get:

$$
\begin{align*}
\sum_{j=1}^{k}\left(F_{j} \Psi_{F_{j}}^{\prime}+\Psi_{F_{j}} F_{j}^{\prime}\right) & =M_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) P_{F, V_{\varepsilon}}^{\prime}+P_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) M_{F, V_{\varepsilon}}^{\prime}+P_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) P_{F, V_{\varepsilon}}^{\prime} \\
& =\left(\Psi_{y}-\Psi_{\varepsilon}\right)-M_{F, V_{\varepsilon}}\left(\Psi_{y}-\Psi_{\varepsilon}\right) M_{F, V_{\varepsilon}}^{\prime} \tag{B.5}
\end{align*}
$$

Then, Equation (B.4) with (B.5) yields Equation (5).
Proof of Proposition 3: We use the following properties of the Hadamard product: $\left(a b^{\prime}\right) \odot$ $\left(c d^{\prime}\right)=(a \odot c)(b \odot d)^{\prime}, \operatorname{diag}\left(a b^{\prime}\right)=a \odot b, a^{\prime} \Delta b=(a \odot b)^{\prime} \operatorname{diag}(\Delta)$, and $(\Delta a) \odot b=a \odot(\Delta b)=$ $\Delta(a \odot b)$ for conformable vectors $a, b, c, d$ and diagonal matrix $\Delta$. Moreover, we deploy the following facts about the vech operator: $\operatorname{diag}\left(G A G^{\prime}\right)=\sqrt{2} \boldsymbol{X}^{\prime} \operatorname{vech}(A)$ and $\operatorname{vech}\left(G^{\prime} \Delta G\right)=$ $\frac{1}{\sqrt{2}} \boldsymbol{X} \operatorname{diag}(\Delta)$ for $(T-k) \times(T-k)$ symmetric matrix $A$ and diagonal $T \times T$ matrix $\Delta$.
(a) With $G=\left[\begin{array}{lllll}g_{1} & : & \cdots & : & g_{T-k}\end{array}\right]$, we have $M_{F, V_{\varepsilon}}=G G^{\prime} V_{\varepsilon}^{-1}=\sum_{j=1}^{T-k} g_{j}\left(V_{\varepsilon}^{-1} g_{j}\right)^{\prime}$. Then, we get the Hadamard product $M_{F, V_{\varepsilon}}^{\odot}=\sum_{i, j=1}^{T-k}\left[g_{i}\left(V_{\varepsilon}^{-1} g_{i}\right)^{\prime}\right] \odot\left[g_{j}\left(V_{\varepsilon}^{-1} g_{j}\right)^{\prime}\right]=\left[\sum_{i, j=1}^{T-k}\left(g_{i} \odot g_{j}\right)\left(g_{i} \odot g_{j}\right)^{\prime}\right]$ $V_{\varepsilon}^{-2}=\left[\sum_{i=1}^{T-k}\left(g_{i} \odot g_{i}\right)\left(g_{i} \odot g_{i}\right)^{\prime}+2 \sum_{i<j}\left(g_{i} \odot g_{j}\right)\left(g_{i} \odot g_{j}\right)^{\prime}\right] V_{\varepsilon}^{-2}=\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right) V_{\varepsilon}^{-2}$.
(b) From part (a) and Equation (7), $\Psi_{\varepsilon}=\operatorname{diag}\left(V_{\varepsilon}^{2}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \operatorname{diag}\left(M_{F, V_{\varepsilon}} Z_{n} M_{F, V_{\varepsilon}}^{\prime}\right)\right)$, with $\operatorname{diag}\left(M_{F, V_{\varepsilon}} Z_{n} M_{F, V_{\varepsilon}}^{\prime}\right)=\operatorname{diag}\left(G G^{\prime} V_{\varepsilon}^{-1} Z_{n} V_{\varepsilon}^{-1} G G^{\prime}\right)=\sqrt{2} \boldsymbol{X}^{\prime} \operatorname{vech}\left(G^{\prime} V_{\varepsilon}^{-1} Z_{n} V_{\varepsilon}^{-1} G\right)$.
(c) We use $\bar{Z}_{n}{ }^{*}=Z_{n}^{*}-G^{\prime} V_{\varepsilon}^{-1} \mathcal{T}_{F, V_{\varepsilon}}\left(Z_{n}\right) V_{\varepsilon}^{-1} G$ in vectorized form. From part (b), we have $\operatorname{diag}\left(\mathcal{T}_{F, V_{\varepsilon}}\left(Z_{n}\right)\right)=\sqrt{2} V_{\varepsilon}^{2}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime} \operatorname{vech}\left(Z_{n}^{*}\right)$. Moreover, $\operatorname{vech}\left(G^{\prime} V_{\varepsilon}^{-1} \mathcal{T}_{F, V_{\varepsilon}}\left(Z_{n}\right) V_{\varepsilon}^{-1} G\right)=$ $\frac{1}{\sqrt{2}} \boldsymbol{X} \operatorname{diag}\left(V_{\varepsilon}^{-1} \mathcal{T}_{F, V_{\varepsilon}}\left(Z_{n}\right) V_{\varepsilon}^{-1}\right)=\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}$ vech $\left(Z_{n}^{*}\right)$. The conclusion follows.

Proof of Proposition 4: (a) We first establish asymptotic normality of $\mathscr{Z}_{n}:=V_{\varepsilon}^{-1 / 2} Z_{n} V_{\varepsilon}^{-1 / 2}$.
Lemma 1 (a) Under Assumptions 1-2, A.2, A. 4 (a)-(b), we have $\Omega_{n}^{-1 / 2}$ vech $\left(\mathscr{L}_{n}\right) \Rightarrow N\left(0, I_{\left.\frac{T(T+1)}{2}\right)}\right)$ as $n \rightarrow \infty$ and $T$ is fixed, where $\Omega_{n}=D_{n}+\kappa_{n} I_{\frac{T(T+1)}{2}}$, and $\kappa_{n}=\frac{1}{n} \sum_{m=1}^{J_{n}}\left(\sum_{i \neq j \in I_{m}} \sigma_{i j}^{2}\right)^{2}$. If additionally Assumption A. 4 (c) holds, then vech $\left(\mathscr{Z}_{n}\right) \Rightarrow N(0, \Omega)$, with $\Omega:=D+\kappa I_{\frac{T(T+1)}{2}}$.

The proof of the next Lemma on vech $\left(Z_{n}^{*}\right)$ being a linear transformation of $\operatorname{vech}\left(\mathscr{Z}_{n}\right)$ uses $\operatorname{vec}(S)=A_{m} \operatorname{vech}(S)$ for any symmetric $m \times m$ matrix $S$, where $A_{m}$ is the $m^{2} \times \frac{1}{2} m(m+1)$ duplication matrix (Magnus, Neudecker (2007)) suited to our definition of the half-vectorization operator vech and given by $A_{m}=\left[\sqrt{2}\left(e_{1} \otimes e_{1}\right): \cdots: \sqrt{2}\left(e_{m} \otimes e_{m}\right):\left\{e_{i} \otimes e_{j}+e_{j} \otimes e_{i}\right\}_{i<j}\right]$, with $e_{i}$ being the $i$ th unit vector in dimension $m$.

Lemma 2 Under Assumption 1, we have vech $\left(Z_{n}^{*}\right)=\operatorname{vech}\left(Q^{\prime} \mathscr{Z}_{n} Q\right)=\boldsymbol{R}^{\prime} \operatorname{vech}\left(\mathscr{Z}_{n}\right)$, where $\boldsymbol{R}:=$ $\frac{1}{2} A_{T}^{\prime}(Q \otimes Q) A_{T-k}$ is a $\frac{1}{2} T(T+1) \times p$ matrix with orthonormal columns, and $Q:=V_{\varepsilon}^{-1 / 2} G$.

From Proposition 3 (c) and Lemma 2, we get $\operatorname{vech}\left(\bar{Z}_{n}^{*}\right)=\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right) \boldsymbol{R}^{\prime} \operatorname{vech}\left(\mathscr{Z}_{n}\right)$. Then, Lemma 1 yields part (a) with $\Omega_{\bar{z}^{*}}=\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right) \boldsymbol{R}^{\prime} \Omega \boldsymbol{R}\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right)$.
(b) We have $\hat{z}_{m, n}=\sum_{i \in I_{m}} \hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1}\left[\tilde{y}_{i} \tilde{y}_{i}^{\prime}-\mathcal{T}_{\hat{F}, \hat{v}_{\varepsilon}}\left(\tilde{y}_{i} \tilde{y}_{i}^{\prime}\right)\right] \hat{V}_{\varepsilon}^{-1} \hat{G}$ with $\tilde{y}_{i}=y_{i}-\bar{y}$, because $\hat{\varepsilon}_{i}=$ $M_{\hat{F}, \hat{\varepsilon}_{\varepsilon}} \tilde{y}_{i}$ and $\hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1} M_{\hat{F}, \hat{V}_{\varepsilon}}=\hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1}$. Using $\tilde{y}_{i} \tilde{y}_{i}^{\prime}=\tilde{\varepsilon}_{i} \tilde{\varepsilon}_{i}^{\prime}+F \beta_{i} \beta_{i}^{\prime} F^{\prime}+F \beta_{i} \tilde{\varepsilon}_{i}^{\prime}+\tilde{\varepsilon}_{i} \beta_{i}^{\prime} F^{\prime}$, we get $\hat{z}_{m, n}=$ $\sum_{i \in I_{m}} \hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1}\left[\tilde{\varepsilon}_{i} \tilde{\varepsilon}_{i}^{\prime}-\mathcal{T}_{\hat{F}, \hat{V}_{\varepsilon}}\left(\tilde{\varepsilon}_{i} \tilde{\varepsilon}_{i}^{\prime}\right)\right] \hat{V}_{\varepsilon}^{-1} \hat{G}+\sum_{i \in I_{m}} \hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1}\left[F \beta_{i} \beta_{i}^{\prime} F^{\prime}-\mathcal{T}_{\hat{F}, \hat{V}_{\varepsilon}}\left(F \beta_{i} \beta_{i}^{\prime} F^{\prime}\right)\right] \hat{V}_{\varepsilon}^{-1} \hat{G}+$ $\sum_{i \in I_{m}} \hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1}\left[F \beta_{i} \tilde{\varepsilon}_{i}^{\prime}+\tilde{\varepsilon}_{i} \beta_{i}^{\prime} F^{\prime}-\mathcal{T}_{\hat{F}, \hat{V}_{\varepsilon}}\left(F \beta_{i} \tilde{\varepsilon}_{i}^{\prime}+\tilde{\varepsilon}_{i} \beta_{i}^{\prime} F^{\prime}\right)\right] \hat{V}_{\varepsilon}^{-1} \hat{G}=: \tilde{z}_{m, n}+z_{m, n, 1}+z_{m, n, 2}$, where $\tilde{\varepsilon}_{i}=\varepsilon_{i}-\bar{\varepsilon}$. Then, we can decompose $\hat{\Omega}_{\bar{Z}^{*}}$ into a sum of a leading term and other terms, which are asymptotically negligible, so that $\hat{\Omega}_{\bar{Z}^{*}}=\tilde{\Omega}_{\bar{Z}^{*}}+o_{p}(1)$, with $\tilde{\Omega}_{\bar{Z}^{*}}=\frac{1}{n} \sum_{m=1}^{J_{n}} \operatorname{vech}\left(\bar{z}_{m, n}\right) \operatorname{vech}\left(\bar{z}_{m, n}\right)^{\prime}$, with $\bar{z}_{m, n}$ defined as $\tilde{z}_{m, n}$ after replacing $\tilde{\varepsilon}_{i}$ with $\varepsilon_{i}$. Let us now show that $\tilde{\Omega}_{\bar{Z}^{*}}=\Omega_{\bar{Z}^{*}}+o_{p}(1)$ up to pre- and post-multiplication by a rotation matrix and its inverse. First, we note that
$\operatorname{vech}\left(\hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1}\left[\varepsilon_{i} \varepsilon_{i}^{\prime}-\mathcal{T}_{\hat{F}, \hat{V}_{\varepsilon}}\left(\varepsilon_{i} \varepsilon_{i}^{\prime}\right)\right] \hat{V}_{\varepsilon}^{-1} \hat{G}\right)=\operatorname{vech}\left(\hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1}\left[\varepsilon_{i} \varepsilon_{i}^{\prime}-\sigma_{i i} V_{\varepsilon}-\mathcal{T}_{\hat{F}, \hat{V}_{\varepsilon}}\left(\varepsilon_{i} \varepsilon_{i}^{\prime}-\sigma_{i i} V_{\varepsilon}\right)\right]\right.$ $\hat{V}_{\varepsilon}^{-1} \hat{G}$ ) because $\mathcal{T}_{\hat{F}, \hat{V}_{\varepsilon}}(\cdot)$ is the identity transformation for diagonal matrices. Moreover, we have:

$$
\begin{align*}
& \text { vech }\left(\hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1}\left[\varepsilon_{i} \varepsilon_{i}^{\prime}-\sigma_{i i} V_{\varepsilon}-\mathcal{T}_{\hat{F}, \hat{V}_{\varepsilon}}\left(\varepsilon_{i} \varepsilon_{i}^{\prime}-\sigma_{i i} V_{\varepsilon}\right)\right] \hat{V}_{\varepsilon}^{-1} \hat{G}\right)=M_{\hat{\mathbf{X}}} \operatorname{vech}\left(\hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1}\left(\varepsilon_{i} \varepsilon_{i}^{\prime}-\sigma_{i i} V_{\varepsilon}\right) \hat{V}_{\varepsilon}^{-1} \hat{G}\right) \\
& \quad=M_{\hat{\boldsymbol{X}}} \text { vech }\left(\tilde{Q}^{\prime}\left(e_{i} e_{i}^{\prime}-\sigma_{i i} I_{T}\right) \tilde{Q}\right)=M_{\hat{\boldsymbol{X}}} \hat{\boldsymbol{R}}^{\prime} \text { vech }\left(e_{i} e_{i}^{\prime}-\sigma_{i i} I_{T}\right), \tag{B.6}
\end{align*}
$$

where $M_{\hat{\boldsymbol{X}}}:=I_{p}-\hat{\boldsymbol{X}}\left(\hat{\boldsymbol{X}}^{\prime} \hat{\boldsymbol{X}}\right)^{-1} \hat{\boldsymbol{X}}^{\prime}$ with $\hat{\boldsymbol{X}}$ defined as $\boldsymbol{X}$ by replacing $G$ with $\hat{G}$, we define $e_{i}=V_{\varepsilon}^{-1 / 2} \varepsilon_{i}$, and $\hat{\boldsymbol{R}}:=\frac{1}{2} A_{T}^{\prime}(\tilde{Q} \otimes \tilde{Q}) A_{T-k}$ with $\tilde{Q}=V_{\varepsilon}^{1 / 2} \hat{V}_{\varepsilon}^{-1} \hat{G}$. The first equality in (B.6) uses an argument similar to Proposition 3 (c), and the third equality is similar to Lemma 2. We get $\operatorname{vech}\left(\bar{z}_{m, n}\right)=M_{\hat{\boldsymbol{X}}} \hat{\boldsymbol{R}}^{\prime} \operatorname{vech}\left(\zeta_{m, n}\right)$, where $\zeta_{m, n}:=\sum_{i \in I_{m}}\left(e_{i} e_{i}^{\prime}-\sigma_{i i} I_{T}\right)$. Besides, vech $\left(\mathscr{Z}_{n}\right)=$ $\frac{1}{\sqrt{n}} \sum_{m=1}^{J_{n}} \operatorname{vech}\left(\zeta_{m, n}\right)$. Then, $\tilde{\Omega}_{\bar{Z}^{*}}=M_{\hat{\boldsymbol{X}}} \hat{\boldsymbol{R}}^{\prime} \tilde{\Omega}_{n} \hat{\boldsymbol{R}} M_{\hat{\boldsymbol{X}}}$ for $\tilde{\Omega}_{n}:=\frac{1}{n} \sum_{m=1}^{J_{n}} \operatorname{vech}\left(\zeta_{m, n}\right) \operatorname{vech}\left(\zeta_{m, n}\right)^{\prime}$. Further, we have $E\left[\tilde{\Omega}_{n}\right]=V\left[\operatorname{vech}\left(\mathscr{Z}_{n}\right)\right]=\Omega_{n}$. Moreover, $\tilde{\Omega}_{n}-E\left[\tilde{\Omega}_{n}\right]=o_{p}(1)$, by using $\operatorname{vec}\left(\tilde{\Omega}_{n}\right)=\frac{1}{n} \sum_{m=1}^{J_{n}} \operatorname{vech}\left(\zeta_{m, n}\right) \otimes \operatorname{vech}\left(\zeta_{m, n}\right)$ and $\left\|V\left[\operatorname{vec}\left(\tilde{\Omega}_{n}\right)\right]\right\| \leq C \frac{1}{n^{2}} \sum_{m=1}^{J_{n}} E\left[\left\|\operatorname{vech}\left(\zeta_{m, n}\right)\right\|^{4}\right]=$ $o(1)$, where the latter bound is shown in the proof of Lemma 1 using Assumption 2 (d). Additionally, by Assumption A.4, we have $\Omega_{n}=\Omega+o(1)$. Thus, $\tilde{\Omega}_{n}=\Omega+o_{p}(1)$. Now, we use consistency of the FA estimates and $\hat{G} \hat{O}=G+o_{p}(1)$ for a (possibly data-dependent) $T-k$ dimensional orthogonal matrix $\hat{O}$. Then, by Proposition 10 (e) in Appendix D.5, we have $\hat{\boldsymbol{R}} M_{\hat{\boldsymbol{X}}} \hat{\mathscr{R}}^{-1}=$ $\hat{\boldsymbol{R}}\left(I_{p}-\hat{\boldsymbol{X}}\left(\hat{\boldsymbol{X}}^{\prime} \hat{\boldsymbol{X}}\right)^{-1} \hat{\boldsymbol{X}}^{\prime}\right) \mathscr{R}(\hat{O})^{-1}=\boldsymbol{R}\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right)+o_{p}(1)$, for a $p$ dimensional orthogonal matrix $\hat{\mathscr{R}} \equiv \mathscr{R}(\hat{O})$. We conclude that $\hat{\mathscr{R}} \tilde{\Omega}_{\bar{Z}^{*}} \hat{\mathscr{R}}^{-1}$ is a consistent estimator of $\Omega_{\bar{Z}^{*}}$ as $n \rightarrow \infty$ and $T$ is fixed, which yields part (b).
(c) Under $H_{1}(k)$ and Assumption A. $5(\mathrm{a})$, we have $\hat{F} \xrightarrow{p} F^{*}$ and $\hat{V}_{\varepsilon} \xrightarrow{p} V_{\varepsilon}^{*}$. Then, $\hat{S} \xrightarrow{p} S^{*}$ with $S^{*}=\left(V_{\varepsilon}^{*}\right)^{-1 / 2} M_{F^{*}, V_{\varepsilon}^{*}}\left(V_{y}-V_{\varepsilon}^{*}\right) M_{F^{*}, V_{\varepsilon}^{*}}^{*}\left(V_{\varepsilon}^{*}\right)^{-1 / 2} \neq 0$. Indeed, if $S^{*}$ were the null matrix, then we would have $M_{F^{*}, V_{\varepsilon}^{*}}\left(V_{y}-V_{\varepsilon}^{*}\right) M_{F^{*}, V_{\varepsilon}^{*}}^{\prime}=0$, which implies $V_{y}-V_{\varepsilon}^{*}=F^{*} A\left(F^{*}\right)^{\prime}$ for a symmetric matrix $A$, in contradiction with Assumption A. 5 (b). Thus, $n\|\hat{S}\|^{2} \geq C n$, w.p.a. 1, for a constant $C>0$. Moreover, using $\operatorname{vech}\left(\hat{z}_{m, n}\right)=\left(I_{p}-\hat{\boldsymbol{X}}\left(\hat{\boldsymbol{X}}^{\prime} \hat{\boldsymbol{X}}\right)^{-1} \hat{\boldsymbol{X}}^{\prime}\right) \operatorname{vech}\left(\hat{G}^{\prime} \hat{V}_{\varepsilon}^{-1}\left(\sum_{i \in I_{m}} \tilde{y}_{i} \tilde{y}_{i}^{\prime}\right) \hat{V}_{\varepsilon}^{-1} \hat{G}\right)$ and the conditions on $\Theta$, we get $\left\|\operatorname{vech}\left(\hat{z}_{m, n}\right)\right\| \leq C \sum_{i \in I_{m}}\left\|\tilde{y}_{i}\right\|^{2}$. Then, from Assumptions A. 1 and A.2, $E\left[\left\|\hat{\Omega}_{\bar{Z}^{*}}\right\|\right] \leq C \frac{1}{n} \sum_{m=1}^{J_{n}} b_{m, n}^{2}=O\left(n \sum_{m=1}^{J_{n}} B_{m, n}^{2}\right)$. Moreover, $\sum_{m=1}^{J_{n}} B_{m, n}^{2}=o(1)$. Indeed, Assumption 2 (d) implies $B_{m, n} \leq c n^{-\frac{\delta}{\delta+1}}$ uniformly in $m$, for any $c>0$ and $n$ large enough, and
hence $\sum_{m=1}^{J_{n}} B_{m, n}^{2}=c n^{-\frac{\delta}{\delta+1}} \sum_{m=1}^{J_{n}} B_{m, n} \leq c$, for any $c>0$ and $n$ large. Part (c) follows.
Proof of Proposition 5: We have $L R(k)=\frac{n}{2}\|\hat{S}\|^{2}+o_{p}(1)=\operatorname{nvech}(\hat{S})^{\prime} \operatorname{vech}(\hat{S})+o_{p}(1)$. Moreover, from the asymptotic expansion (9), we can write $\sqrt{n} v e c h(\hat{S})=A\left(F, V_{\varepsilon}\right) z_{n}^{A D}+o_{p}(1)$, where vector $z_{n}^{A D}$ stacks the $T(T-1) / 2$ above-diagonal elements of matrix $Z_{n}$ and $A\left(F, V_{\varepsilon}\right)$ is a deterministic matrix whose elements only depend on $F, V_{\varepsilon}$. From Conditions (a) and (b) of Proposition 5, and Lemma 1, we have $z_{n}^{A D} \Rightarrow N\left(0, \Omega_{z}\right)$, where the diagonal matrix $\Omega_{z}$ is the same as if the errors were independent normally distributed - up to replacing $q$ with $q+\kappa$.

Proof of Proposition 6: Let us first get the asymptotic expansion of $\hat{V}_{y}=\frac{1}{n} \tilde{Y} \tilde{Y}^{\prime}$. With the drifting DGP $Y=\mu 1_{n}^{\prime}+F \beta^{\prime}+F_{k+1} \beta_{l o c}^{\prime}+\varepsilon$, and using $\bar{\beta}=0, \bar{\beta}_{l o c}=0, \frac{1}{n}\left[\beta: \beta_{l o c}\right]^{\prime}\left[\beta: \beta_{l o c}\right]=I_{k+1}$ and Lemma 6 (a) in Appendix D, we get $\hat{V}_{y}=\tilde{V}_{y}+\frac{1}{\sqrt{n}} \Psi_{y, l o c}+R_{y}$, where $\tilde{V}_{y}=F F^{\prime}+\tilde{V}_{\varepsilon}$,

$$
\begin{equation*}
\Psi_{y, l o c}=c_{k+1} \rho_{k+1} \rho_{k+1}^{\prime}+\frac{1}{\sqrt{n}}\left(\varepsilon \beta F^{\prime}+F \beta^{\prime} \varepsilon^{\prime}\right)+\sqrt{n}\left(\frac{1}{n} \varepsilon \varepsilon^{\prime}-\tilde{V}_{\varepsilon}\right), \tag{B.7}
\end{equation*}
$$

and $R_{y}=\frac{1}{n}\left(\varepsilon \beta_{l o c} F_{k+1}^{\prime}+F_{k+1} \beta_{l o c}^{\prime} \varepsilon^{\prime}\right)+\left[F_{k+1} F_{k+1}^{\prime}-n^{-1 / 2} c_{k+1} \rho_{k+1} \rho_{k+1}^{\prime}\right]+o_{p}\left(\frac{1}{\sqrt{n}}\right)$. Using $F_{k+1}=$ $\sqrt{\gamma_{k+1}} \rho_{k+1}$ and $\sqrt{n} \gamma_{k+1}=c_{k+1}+o(1)$, we get $R_{y}=o_{p}(1 / \sqrt{n})$. We use Equation (5) with $\Psi_{y, l o c}$ given in (B.7) instead of $\Psi_{y}$, and get $\Psi_{\varepsilon}=\mathcal{T}_{F, V_{\varepsilon}}\left(Z_{n, l o c}\right)$, where $Z_{n, l o c}:=\sqrt{n}\left(\frac{1}{n} \varepsilon \varepsilon^{\prime}-\tilde{V}_{\varepsilon}\right)+$ $c_{k+1} \rho_{k+1} \rho_{k+1}^{\prime}$ and $\operatorname{diag}\left(\Psi_{\varepsilon}\right)=\sqrt{2} V_{\varepsilon}^{2}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime} \operatorname{vech}\left(Z_{n, l o c}^{*}\right)$. Then, as in Equation (8), $\sqrt{n} \hat{S}=$ $V_{\varepsilon}^{-1 / 2} M_{F, V_{\varepsilon}}\left(Z_{n, l o c}-\mathcal{T}_{F, V_{\varepsilon}}\left(Z_{n, l o c}\right)\right) M_{F, V_{\varepsilon}}^{\prime} V_{\varepsilon}^{-1 / 2}+o_{p}(1)=V_{\varepsilon}^{-1 / 2} G\left(\bar{Z}_{n}^{*}+\Delta\right) G^{\prime} V_{\varepsilon}^{-1 / 2}+o_{p}(1)$, where $\Delta=c_{k+1} G^{\prime} V_{\varepsilon}^{-1}\left(\rho_{k+1} \rho_{k+1}^{\prime}-\mathcal{T}_{F, V_{\varepsilon}}\left(\rho_{k+1} \rho_{k+1}^{\prime}\right)\right) V_{\varepsilon}^{-1} G$. As in Proposition 3 (c), we have $\operatorname{vech}(\Delta)=c_{k+1}\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right) \operatorname{vech}\left(G^{\prime} V_{\varepsilon}^{-1} \rho_{k+1} \rho_{k+1}^{\prime} V_{\varepsilon}^{-1} G\right)=c_{k+1}\left(I_{p}-\boldsymbol{X}\left(\boldsymbol{X}^{\prime} \boldsymbol{X}\right)^{-1} \boldsymbol{X}^{\prime}\right)$ $\operatorname{vech}\left(\xi_{k+1} \xi_{k+1}^{\prime}\right)$ since $\rho_{k+1}=G \xi_{k+1}$. From the proof of Proposition 4 (a), $\bar{Z}_{n}^{*} \Rightarrow \bar{Z}^{*}$ as $n \rightarrow \infty$, and Part (a) follows. Part (b) is a consequence of Part (a) and the Continuous Mapping Theorem.

Proof of Proposition 7: The proof of part (a) is in three steps. (i) The testing problem asymptotically simplifies to the null hypothesis $H_{0}: \lambda_{1}=\ldots=\lambda_{d f}=0$ vs. the alternative hypothesis $H_{1}: \exists \lambda_{j}>0, j=1, \ldots, d f$. Let us define $\boldsymbol{\lambda}_{0}=(0, \ldots, 0)^{\prime}$ for the null hypothesis and pick a given vector $\boldsymbol{\lambda}_{1}=\left(\lambda_{1}, \ldots, \lambda_{d f}\right)^{\prime}$ in the alternative hypothesis, and consider the test of $\boldsymbol{\lambda}_{0}$ versus $\boldsymbol{\lambda}_{1}$ (simple hypothesis). By Neyman-Pearson Lemma, the most powerful test for $\boldsymbol{\lambda}_{0}$ versus $\boldsymbol{\lambda}_{1}$ rejects the null hypothesis when $f\left(z ; \lambda_{1}, \ldots, \lambda_{d f}\right) / f(z ; 0, \ldots, 0)$ is large, i.e., the test function is
$\phi(z)=1\left\{\frac{f\left(z ; \lambda_{1}, \ldots, \lambda_{d f}\right)}{f(z ; 0, \ldots, 0)} \geq C\right\}$ for a constant $C>0$ set to ensure the correct asymptotic size.
(ii) Let us now show that the density ratio $\frac{f\left(z ; \lambda_{1}, \ldots, \lambda_{d f}\right)}{f(z ; 0, \ldots, 0)}$ is an increasing function of $z$. To show this, we can rely on an expansion of the density of $\sum_{j=1}^{d f} \mu_{j} \chi^{2}\left(1, \lambda_{j}^{2}\right)$ in terms of central chi-square densities (Kotz, Johnson, and Boyd (1967) Equations (144) and (151)):

$$
\begin{equation*}
f\left(z ; \lambda_{1}, \ldots, \lambda_{d f}\right)=\sum_{k=0}^{\infty} \bar{c}_{k}\left(\lambda_{1}, \ldots, \lambda_{d f}\right) g(z ; d f+2 k, 0) \tag{B.8}
\end{equation*}
$$

where the coefficients $\bar{c}_{k}\left(\lambda_{1}, \ldots, \lambda_{d f}\right)=A e^{-\sum_{j=1}^{d f} \lambda_{j}^{2} / 2} E\left[Q\left(\lambda_{1}, \ldots, \lambda_{d f}\right)^{k}\right] / k$ ! involve moments of the quadratic form $Q\left(\lambda_{1}, \ldots, \lambda_{d f}\right)=(1 / 2) \sum_{j=1}^{d f}\left(\nu_{j}^{1 / 2} X_{j}+\lambda_{j}\left(1-\nu_{j}\right)^{1 / 2}\right)^{2}$ of the mutually independent variables $X_{j} \sim N(0,1), A=\prod_{j=1}^{d f} \mu_{j}^{-1 / 2}$, and $\nu_{j}=1-\frac{1}{\mu_{j}} \min _{\ell} \mu_{\ell}$. Without loss of generality for checking the monotonicity, we have rescaled the density so that $\min _{j} \mu_{j}=$ 1. Then, from (B.8), we get the ratio: $\frac{f\left(z ; \lambda_{1}, \ldots, \lambda_{d f}\right)}{f(z ; 0, \ldots, 0)}=\frac{\sum_{k=0}^{\infty} \bar{c}_{k}\left(\lambda_{1}, \ldots, \lambda_{d f}\right) g(z ; d f+2 k, 0)}{\sum_{k=0}^{\infty} \bar{c}_{k}(0, \ldots, 0) g(z ; d f+2 k, 0)}$. By dividing both the numerator and the denominator by the central chi-square density $g(z ; d f, 0)$, we get $\frac{f\left(z ; \lambda_{1}, \ldots, \lambda_{d f}\right)}{f(z ; 0, \ldots, 0)}=e^{-\sum_{j=1}^{d f} \lambda_{j}^{2} / 2 \frac{\sum_{k=0}^{\infty} c_{k}\left(\lambda_{1}, \ldots, \lambda_{d f}\right) \psi_{k}(z)}{\sum_{k=0}^{\infty} c_{k}(0, \ldots, 0) \psi_{k}(z)}}=: e^{-\sum_{j=1}^{d f} \lambda_{j}^{2} / 2} \Psi\left(z ; \lambda_{1}, \ldots, \lambda_{d f}\right)$, where $\psi_{k}(z):=$ $g(z ; d f+2 k, 0) / g(z ; d f, 0)=\frac{\Gamma\left(\frac{d f}{2}\right)}{2^{k} \Gamma\left(\frac{d f}{2}+k\right)} z^{k}$ is the ratio of central chi-square distributions with $d f+2 k$ and $d f$ degrees of freedom, and $c_{k}\left(\lambda_{1}, \ldots, \lambda_{d f}\right)=E\left[Q\left(\lambda_{1}, \ldots, \lambda_{d f}\right)^{k}\right] / k!$. We use the short notation $c_{k}(\lambda):=c_{k}\left(\lambda_{1}, \ldots, \lambda_{d f}\right)$ and $c_{k}(0):=c_{k}(0, \ldots, 0)$. The factor $e^{-\sum_{j=1}^{d f} \lambda_{j}^{2} / 2}$ does not impact on the monotonicity of the density ratio. We take the derivative of $\Psi\left(z ; \lambda_{1}, \ldots, \lambda_{d f}\right)$ with respect to argument $z$ and get $\partial_{z} \Psi\left(z ; \lambda_{1}, \ldots, \lambda_{d f}\right)=\frac{\left(\sum_{k=1}^{\infty} c_{k}(\lambda) \psi_{k}^{\prime}(z)\right)\left(1+\sum_{k=1}^{\infty} c_{k}(0) \psi_{k}(z)\right)}{\left(\sum_{k=0}^{\infty} c_{k}(0) \psi_{k}(z)\right)^{2}}-$ $\frac{\left(1+\sum_{k=1}^{\infty} c_{k}(\lambda) \psi_{k}(z)\right)\left(\sum_{k=1}^{\infty} c_{k}(0) \psi_{k}^{\prime}(z)\right)}{\left(\sum_{k=0}^{\infty} c_{k}(0) \psi_{k}(z)\right)^{2}}$. The sign is given by the difference of the numerators, which is $\sum_{k=1}^{\infty}\left[c_{k}(\lambda)-c_{k}(0)\right] \psi_{k}^{\prime}(z)+\sum_{k, l=1, k \neq l}^{\infty} c_{k}(\lambda) c_{l}(0)\left[\psi_{k}^{\prime}(z) \psi_{l}(z)-\psi_{k}(z) \psi_{l}^{\prime}(z)\right]=\sum_{k=1}^{\infty}\left[c_{k}(\lambda)-\right.$ $\left.c_{k}(0)\right] \psi_{k}^{\prime}(z)+\sum_{k, l=1, k>l}^{\infty}\left[c_{k}(\lambda) c_{l}(0)-c_{l}(\lambda) c_{k}(0)\right]\left[\psi_{k}^{\prime}(z) \psi_{l}(z)-\psi_{k}(z) \psi_{l}^{\prime}(z)\right]$. We use $\psi_{k}^{\prime}(z)=$ $\frac{\Gamma\left(\frac{d}{2}\right) k}{2^{k} \Gamma\left(\frac{d}{2}+k\right)} z^{k-1}$ and $\psi_{k}^{\prime}(z) \psi_{l}(z)-\psi_{k}(z) \psi_{l}^{\prime}(z)=(k-l) \frac{\Gamma\left(\frac{d}{2}\right)^{2}}{2^{k+l} \Gamma\left(\frac{d}{2}+k\right) \Gamma\left(\frac{d}{2}+l\right)} z^{k+l-1}$ for $k>l$ and $z \geq 0$. The difference of the numerators in the derivative of the density ratio becomes: $\frac{1}{2} \frac{\Gamma\left(\frac{d}{2}\right)}{\Gamma\left(\frac{d}{2}+1\right)}\left[c_{1}(\lambda)-c_{1}(0)\right]+\frac{1}{2^{2}} \frac{2 \Gamma\left(\frac{d}{2}\right)}{\Gamma\left(\frac{d}{2}+2\right)}\left[c_{2}(\lambda)-c_{2}(0)\right] z+\sum_{m=3}^{\infty} \frac{1}{2^{m}}\left(m \frac{\Gamma\left(\frac{d}{2}\right)}{\Gamma\left(\frac{d}{2}+m\right)}\left[c_{m}(\lambda)-c_{m}(0)\right]\right.$ $\left.+\sum_{k>l \geq 1, k+l=m} \frac{(k-l) \Gamma\left(\frac{d}{2}\right)^{2}}{\Gamma\left(\frac{d}{2}+k\right) \Gamma\left(\frac{d}{2}+l\right)}\left[c_{k}(\lambda) c_{l}(0)-c_{l}(\lambda) c_{k}(0)\right]\right) z^{m-1}=\sum_{m=1}^{\infty} \frac{1}{2^{m}} \kappa_{m} z^{m-1}$, with $\kappa_{m}:=$
$\sum_{k>l \geq 0, k+l=m}(k-l) \frac{\Gamma\left(\frac{d}{2}\right)^{2}}{\Gamma\left(\frac{d}{2}+k\right) \Gamma\left(\frac{d}{2}+l\right)}\left[c_{k}(\lambda) c_{l}(0)-c_{l}(\lambda) c_{k}(0)\right]$. A direct calculation shows that $\kappa_{1}, \kappa_{2} \geq$ 0 . Hence, a sufficient condition for monotonicity of the density ratio is $\kappa_{m} \geq 0$, for all $m \geq 3$, i.e., Inequalities (12). Thus, the test rejects for large values of the argument, i.e., $\phi(z)=1\{z \geq \bar{C}\}$, where the constant $\bar{C}$ is determined by fixing the asymptotic size under the null hypothesis.
(iii) Since the test function $\phi$ does not depend on $\boldsymbol{\lambda}_{1}$, it is AUMPI in the class of hypothesis tests based on the LR statistic (or the squared norm statistic). It yields part (a).

Let us now turn to the proof of part (b). From the definition of the $\kappa_{m}$ coefficients written as $\kappa_{m}=\sum_{j>l \geq 0, j+l=m} \frac{(j-l) \Gamma\left(\frac{d f}{2}\right)^{2}}{\Gamma\left(\frac{d f}{2}+j\right) \Gamma\left(\frac{d f}{2}+l\right)} c_{j}(0) c_{l}(0)\left[\frac{c_{j}(\lambda)}{c_{j}(0)}-\frac{c_{l}(\lambda)}{c_{l}(0)}\right]$, it is sufficient to get $\kappa_{m} \geq 0$, for all $m$, that sequence $\frac{c_{j}(\lambda)}{c_{j}(0)}$, for $j=0,1, \ldots$, is increasing. To prove that, we link the coefficients $c_{j}(\lambda)$ to the complete exponential Bell's polynomials (Bell (1934)) and establish the following recurrence.

Lemma 3 We have $c_{l+1}(\lambda)=\frac{1}{l+1} \sum_{i=0}^{l}\left(\frac{1}{2} \sum_{j=1}^{d f} \nu_{j}^{i}\left[\nu_{j}+(i+1)\left(1-\nu_{j}\right) \lambda_{j}^{2}\right]\right) c_{l-i}(\lambda)$, for $l \geq 0$.
We use $\frac{c_{l}(\lambda)}{c_{l}(0)}=\frac{\tilde{c}_{l}(\lambda)}{\gamma_{l}}$, where we obtain the sequences $\gamma_{l}:=c_{l}(0) \nu_{d f}^{-l}$ and $\tilde{c}_{l}(\lambda):=c_{l}(\lambda) \nu_{d f}^{-l}$ by standardization with $\nu_{d f}^{-l}$. From Lemma 3, we have $\gamma_{l+1}=\frac{1}{l+1} \sum_{i=0}^{l} \frac{1}{2}\left(1+\sum_{j=2}^{d f-1} \rho_{j}^{i+1}\right) \gamma_{l-i}$ with $\gamma_{0}=1$, and $\tilde{c}_{l+1}(\lambda)=\frac{1}{l+1} \sum_{i=0}^{l}\left(\frac{1}{2} \sum_{j=1}^{d f} \rho_{j}^{i}\left[\rho_{j}+\frac{i+1}{\nu_{d f}}\left(1-\nu_{j}\right) \lambda_{j}^{2}\right]\right) \tilde{c}_{l-i}(\lambda)$ with $\tilde{c}_{0}(\lambda)=1$ (note that $\rho_{1}=0$ and $\rho_{d f}=1$ ). To prove that sequence $\frac{\tilde{c}_{l}(\lambda)}{\gamma_{l}}$ is increasing, the next lemma provides a sufficient condition from "separation" of the coefficients that define the recursive relations.

Lemma 4 (Separation Lemma) Let $\left(a_{i}\right)$ be a real sequence, and let $b_{i}=\frac{1}{2}\left(1+\sum_{j=2}^{d f-1} \rho_{j}^{i}\right)$, for $i \geq 1$, where $0 \leq \rho_{j} \leq 1$. Let sequences $\left(g_{l}\right)$ and $\left(c_{l}\right)$ be defined recursively by $g_{l+1}=$ $\frac{1}{l}\left(b_{1} g_{l}+b_{2} g_{l-1}+\ldots+b_{l}\right)$ and $c_{l+1}=\frac{1}{l}\left(a_{1} c_{l}+a_{2} c_{l-1}+\ldots+a_{l}\right)$, with $g_{1}=c_{1}=1$. Suppose that $a_{i} \geq \max \left\{\frac{d f-1}{2}, 1\right\}$, for all $i$ (separation condition). Then, sequence $\left(\frac{c_{l}}{g_{l}}\right)$ is increasing.

We apply Lemma 4 to sequences $\tilde{c}_{l}(\lambda)$ and $\gamma_{l}$. We detail the case $d f \geq 3$ (for $d f=2$ the analysis is simpler). The separation condition $\frac{1}{2} \sum_{j=1}^{d f} \rho_{j}^{i}\left[\rho_{j}+\frac{i+1}{\nu_{d f}}\left(1-\nu_{j}\right) \lambda_{j}^{2}\right] \geq \frac{d f-1}{2}$, for $i=0$, yields $\lambda_{1}^{2}+\sum_{j=2}^{d f}\left(1-\nu_{j}\right) \lambda_{j}^{2} \geq \nu_{d f}\left(d f-2-\sum_{j=2}^{d f-1} \rho_{j}\right)$, and, for $i \geq 1$, it yields $\sum_{j=2}^{d f-1} \rho_{j}^{i}(1-$ $\left.\nu_{j}\right) \lambda_{j}^{2}+\left(1-\nu_{d f}\right) \lambda_{d f}^{2} \geq \frac{\nu_{d f}}{i+1}\left(d f-2-\sum_{j=2}^{d f-1} \rho_{j}^{i+1}\right)$. Inequalities (13) follow.

Figure 1: The upper panel displays the p-values for the statistic $L R(k)$ for the subperiods from January 1963 to December 2021, stopping at the smallest $k$ such that $H_{0}(k)$ is not rejected at level $\alpha_{n}=10 / n_{\max }$. If no such $k$ is found then p-values are displayed up to $k_{\max }$. We use rolling windows of $T=20$ months moving forward by 12 months each time. The first bar of p-values covers the whole 20 months. Other bars cover the last 12 months of the 20 months subperiod. We flag bear market phases with grey shaded vertical bars. The five lower panels display $\left(\hat{\hat{V}}_{y}\right)^{1 / 2}$ for total cross-sectional volatility, $\left(\overline{\hat{F}^{\prime}} \hat{F}\right)^{1 / 2}$ for systematic volatility, $\left(\overline{\hat{V}}_{\varepsilon}\right)^{1 / 2}$ for idiosyncratic volatility, as well as $\hat{R}^{2}$ and $\hat{R}^{2}$ under a single-factor model.


Systematic Volatility



Single-Factor Model $\hat{R}^{2}$


Figure 2: The upper and lower panels display the p-values for the RS and KP statistics for the subperiods from January 1963 to December 2021, for the rank test of the null hypothesis $H_{0, s p}(r)$ that $F^{O}$ has rank $r$ against the alternative hypothesis of rank larger than $r$, for any integer $r \leq k-1$. The empirical matrix $\hat{F}^{O}$ is computed with the time-varying portfolio weights of the Fama-French five-factor model plus momentum. We stop at the smallest $r$ such that $H_{0, s p}(r)$ is not rejected at level $\alpha_{n}=10 / n$. If no such $r$ is found then p-values are displayed up to $k-1$. The red horizontal segments give $\hat{k}-1$, i.e., the estimated number of latent factors obtained from Figure 1 minus 1. We flag bear market phases with grey shaded vertical bars, and use the same rolling windows as in Figure 1.



[^0]:    ${ }^{* 1}$ University of Geneva, ${ }^{2}$ Swiss Finance Institute, ${ }^{3}$ Università della Svizzera italiana. Acknowledgements: We are grateful to A. Onatski for his very insightful discussion of our paper Fortin, Gagliardini, Scaillet (2022) at the 14th Annual SoFiE Conference in Cambridge, which prompted us to exploit factor analysis for estimation in short panels. We thank L. Barras, S. Bonhomme, M. Caner, F. Carlini, I. Chaieb, F. Ghezzi, A. Horenstein, S. Kim, F. Kleibergen, H. Langlois, T. Magnac, S. Ng, E. Ossola, Y. Potiron, F. Trojani, participants at (EC)^2 2022, QFFE 2023, SoFiE 2023, NASM 2023, IPDC 2023, and participants at seminars at UNIMIB, UNIBE, UNIGE, Warwick University.

[^1]:    ${ }^{1}$ Raponi, Robotti and Zaffaroni (2020) develop tests of beta-pricing models and a two-pass methodology to estimate

[^2]:    ${ }^{2}$ Under NA, the intercept term in the asset return model $y_{i}=\mu_{i}+\tilde{F} \tilde{\beta}_{i}+\varepsilon_{i}$ is $\mu_{i}=r_{f}+1_{T} \nu^{\prime} \tilde{\beta}_{i}$, where $r_{f}$ is the $T$-dimensional vector whose entries collect the (possibly time-varying) risk-free rates, $\nu=\left(\nu_{1}, \ldots, \nu_{k}\right)^{\prime}$ is a $k$ dimensional vector of parameters, and $1_{T}$ is a $T$-dimensional vector of ones (see e.g. Gagliardini, Ossola and Scaillet (2016)). We can absorb term $1_{T} \nu^{\prime} \tilde{\beta}_{i}$ into the systematic part to get $y_{i}=r_{f}+F \tilde{\beta}_{i}+\varepsilon_{i}$ with $F=\tilde{F}+1_{T} \nu^{\prime}$. It holds irrespective of the latent factors being tradable or not. If the factors are tradable, we further have $\nu=0$ from the NA restriction. Akin to standard formulation of FA, we recenter the latent effects by subtracting their mean $\tilde{\mu}_{\tilde{\beta}}=\frac{1}{n} \sum_{i=1}^{n} \tilde{\beta}_{i}$, to get model (1) with $\beta_{i}=\tilde{\beta}_{i}-\tilde{\mu}_{\tilde{\beta}}$ and $\mu=r_{f}+F \tilde{\mu}_{\tilde{\beta}}$.
    ${ }^{3}$ When there is a common random component in idiosyncratic volatilities, we have $V_{\varepsilon}=\operatorname{plim}_{n \rightarrow \infty} \frac{1}{n} \varepsilon \varepsilon^{\prime}$ by a suitable version of the Law of Large Number (LLN) conditional on the sigma-field generated by this common component. With fixed $T$, we treat the sample realizations of the common component in idiosyncratic volatilities as unknown time fixed effects (the diagonal elements of matrix $V_{\varepsilon}$ ), which yields time heterogeneous distributions for the errors. It is how the unconditional expectation in Assumption 1 has to be understood.

[^3]:    ${ }^{4}$ Chamberlain (1992) studies semiparametrically efficient estimation in panel models with fixed effects and short $T$ using moment restrictions from instrumental variables. Our approach does not rely on availability of valid instruments.

[^4]:    ${ }^{5}$ The normalization in (FA2) applies for $\hat{\gamma}_{j} \geq 0$, which holds with probability approaching 1. Otherwise, the first-order conditions of the FA estimators hold with $\hat{\gamma}_{j}$ replaced by its positive part; see Anderson (2003) for similar positivity constraints.
    ${ }^{6}$ Integer $d f$ is the number of different elements in data matrix $\hat{V}_{y}$, i.e., $\frac{1}{2} T(T+1)$, plus the number of normalization constraints $\frac{1}{2} k(k-1)$ in equations $F^{\prime} V_{\varepsilon}^{-1} F=\operatorname{diag}$, minus the number of unknown parameters $(k+1) T$ (Anderson (2003)). In SMC Table 3, we list the largest admissible number $k$ of latent factors as a function of $T$ such that $d f \geq 0$.
    ${ }^{7}$ It comes from second-order expansion of the $\log$ function using $\sum_{j=k+1}^{T} \hat{\gamma}_{j}=0$, which is a consequence of Proposition 1 (a) and (c) (see also Anderson (2003)), and $\sqrt{n} \hat{\gamma}_{j}=O_{p}(1)$ for $j=k+1, \ldots, T$, which follows from Propositions 1 and 4.

[^5]:    ${ }^{8}$ The test in Connor and Korajczyk (1993) is built on cross-sectional averages of squared residuals, akin to diagonal terms of $\hat{S}$, but obtained by PCA instead of FA. However, their test statistic involves the difference of such crosssectional averages for two consecutive dates, and relies on error sphericity. Furthermore, from Proposition 1 (a), we note that test statistics based on the spacings between the small eigenvalues of $\hat{V}_{y} \hat{V}_{\varepsilon}^{-1}$ use the non-zero eigenvalues of $\hat{S}$. Such tests rely on the possibility to identify the number of latent factors $k$ from the fact that the $k$ smallest eigenvalues of $V_{y} V_{\varepsilon}^{-1}$ are all equal to 1.

[^6]:    ${ }^{9}$ The remaining eigenvalue is equal to 1 with multiplicity $T-k$. We have $F_{j}=\sqrt{\gamma_{j}} V_{\varepsilon}^{1 / 2} U_{j}$, where the $U_{j}$ are the orthonormal eigenvectors of $V_{\varepsilon}^{-1 / 2} V_{y} V_{\varepsilon}^{-1 / 2}$ for the $k$ largest eigenvalues $1+\gamma_{j}$.

[^7]:    ${ }^{10} \mathrm{We}$ can extend results like (10) to test statistics that are generic functions of the eigenvalues of matrix $\hat{S}$ by using the Weyl's inequalities (see e.g. Bernstein (2009)), and develop test statistics along the lines of FGS.

[^8]:    ${ }^{11}$ This definition of the half-vectorization operator for symmetric matrices differs from the usual one by the ordering of the elements, and the rescaling of the diagonal elements. It is more convenient for our purposes (see proof of Lemma 11). For instance, it holds $\frac{1}{2}\|A\|^{2}=\operatorname{vech}(A)^{\prime} \operatorname{vech}(A)$, for a symmetric matrix $A$.

[^9]:    ${ }^{12}$ In Appendix D.4.3, we present a parametric estimator for the asymptotic variance of $\operatorname{vech}\left(Z_{n}\right)$ under additional conditions on the error terms
    ${ }^{13}$ In Proposition 10 in Appendix D.5, we study how $\boldsymbol{X}$ and $\Omega_{\bar{Z}^{*}}$ are transformed under different choices for the rotation of $G$. The eigenvalues $\mu_{j}$ are invariant to such rotation as expected.

[^10]:    ${ }^{14}$ For instance, we can set $\hat{Q}=\tilde{Q}\left(\tilde{Q}^{\prime} \tilde{Q}\right)^{-1 / 2}$, where matrix $\tilde{Q}$ consists of the first $T-k$ columns of $I_{T}-\hat{V}_{\varepsilon}^{-1 / 2} \hat{F}\left(\hat{F}^{\prime} \hat{V}_{\varepsilon}^{-1} \hat{F}\right)^{-1} \hat{F}^{\prime} \hat{V}_{\varepsilon}^{-1 / 2}$, if those columns are linearly independent.

[^11]:    ${ }^{15}$ If the $\sigma_{i i}$ were treated as i.i.d. random effects independent of errors, and we exclude cross-sectional correlation of errors to simplify, we would recover the i.i.d. condition of the data. However, the random $\sigma_{i i}$ would yield a stochastic common factor across time that breaks the condition in Corollary 2 of Anderson and Amemiya (1988).

[^12]:    ${ }^{16}$ We have $\hat{\gamma}_{j}=\hat{\delta}_{j} / \hat{\sigma}^{2}-1$, and $L R(k)=\frac{n}{2 \hat{\sigma}^{4}} \sum_{j=k+1}^{T}\left(\hat{\delta}_{j}^{2}-\hat{\sigma}^{2}\right)^{2}+o_{p}(1)$, i.e., the LR statistic is asymptotically equivalent to the sum of squared deviations of the $T-k$ smallest eigenvalues from their mean. Besides, by similar results, we have that the eigenvalue spacing statistic $\sqrt{n} \mathscr{S}(k):=\hat{\gamma}_{k+1}-\hat{\gamma}_{T}$ corresponds to the statistic considered in FGS divided by $\hat{\sigma}^{2}$, and its asymptotic distribution coincides with that obtained by FGS.

[^13]:    ${ }^{17}$ Here, we do not deal with invariance to data transformations but rather with invariance to parameterization of $G$

[^14]:    and $\boldsymbol{D}$. However, if we consider tests based on the elements of vector $\hat{W}$, this difference is immaterial.

[^15]:    ${ }^{18}$ Inequalities (12) with $d f=1$ are easily proved to hold. In such a case, we can use the asymptotic distribution of a scaled chi-square variable and its MLR property.

[^16]:    ${ }^{19}$ We fix their parameter values $\lambda_{1}=\lambda_{2}=0.2$ for the classification based on the nominal S\&P500 index. Bear periods are close to NBER recessions.

[^17]:    ${ }^{20} \mathrm{We}$ also investigate stability of the factor structure by dividing each window of 20 months into two overlapping subperiods of 16 months (overlap of 12 months) and by estimating canonical correlations between the betas in each subperiod (see SMC ). We find that the fraction of common factors is 1 in $70 \%$ of the windows. The fraction is between 0.8 and 1 in $25 \%$. It is between 0.5 and 0.8 in the remaining periods.

[^18]:    ${ }^{21}$ With fixed $T$, the selection procedure of Zaffaroni (2019), being by construction more conservative than a (multiple) testing procedure (see the discussion on p. 508 of Gagliardini, Scaillet, and Ossola (2019)), yields a smaller number of factors. Imposing cross-sectional independence (resp., Gaussianity and cross-sectional independence) for the LR test gives most of the time an increase by 1 or 2 ( 1 or 3 ). We have an average increase of 3 under sphericity.
    ${ }^{22}$ This choice satisfies the theoretical rule $\log \alpha_{n} / n \rightarrow 0$ given in Pötscher (1983).
    ${ }^{23} \mathrm{We}$ do not plot the whole paths date $t$ by date $t$, but only averages, for readability. If we sum over time instead of averaging the estimated variances, we get a quantity similar to an integrated volatility (see e.g. Barndorff-Nielsen and Shephard (2002), Andersen, Bollerslev, Diebold, and Labys (2003), and references in Aït-Sahalia and Jacod (2014)), and $\hat{R}^{2}$ is the ratio of such quantities.

[^19]:    ${ }^{24}$ Cross-sectional independence (resp. cross-sectional independence and Gaussianity / sphericity) increases estimated systematic risk in average by $0.6 \%$ (resp. $0.6 \% / 1 \%$ ) and decreases estimated idiosyncratic volatility in average by $0.5 \%$ (resp. $0.5 \% / 1.3 \%$ ), so that estimated $R^{2}$ is inflated in average by $4 \%$ (resp. $4 \% / 12.8 \%$ ).
    ${ }^{25}$ As in Campbell et al. (2023), we have also made the estimation on value-weighted returns and we confirm that the results are qualitatively similar.

[^20]:    ${ }^{26}$ If $k^{O}<k$, empirical factors cannot span the latent space by construction. The condition $k^{O} \geq k$ eases discussion but is not needed for the rank tests.
    ${ }^{27}$ Ahn, Horenstein and Wang (2018) use that technology in a fixed- $n$ large- $T$ setting, and find that ranks of beta matrices estimated from either portfolios, or individual stocks, excess returns are often substantially smaller than the (potentially large) number $k^{O}$ of observed factors. The explanation in large economies is that the portfolio beta matrices coincide with $\Phi$, and thus they cannot have a rank above the (potentially small) number $k$ of latent factors.
    ${ }^{28}$ Spanning holds if we can reject $H_{0, s p}(r)$ for any $r<k$.

[^21]:    ${ }^{29} \operatorname{Under} \operatorname{Rank}(F)=k$, we have $\operatorname{Rank}\left(F^{O}\right)=\operatorname{Rank}(\Phi)$. Hence, under $H_{0, s p}(r)$, matrix $\Phi$ has reduced rank $r$.

