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# LOCAL ASYMPTOTIC NORMALITY OF THE SPECTRUM OF HIGH-DIMENSIONAL SPIKED F-RATIOS 

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#### Abstract

We consider two types of spiked multivariate F distributions: a scaled distribution with the scale matrix equal to a rank-k perturbation of the identity, and a distribution with trivial scale, but rank-k noncentrality. The eigenvalues of the rank-k matrix (spikes) parameterize the joint distribution of the eigenvalues of the corresponding F matrix. We show that, for the spikes located above a phase transition threshold, the asymptotic behavior of the log ratio of the joint density of the eigenvalues of the F matrix to their joint density under a local deviation from these values depends only on the $k$ of the largest eigenvalues \{code\} Furthermore, we show that \{code\} are asymptotically jointly normal, and the statistical experiment of observing all the eigenvalues of the F matrix converges in the Le Cam sense to a Gaussian shift experiment that depends on the asymptotic means and variances of \{code\} . In particular, the best statistical inference about sufficiently large spikes in the local asymptotic regime is based on the $k$ of the largest eigenvalues only.


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#### Abstract

We consider two types of spiked multivariate F distributions: a scaled distribution with the scale matrix equal to a rank- $k$ perturbation of the identity, and a distribution with trivial scale, but rank- $k$ non-centrality. The eigenvalues of the rank- $k$ matrix (spikes) parameterize the joint distribution of the eigenvalues of the corresponding F matrix. We show that, for the spikes located above a phase transition threshold, the asymptotic behavior of the log ratio of the joint density of the eigenvalues of the F matrix to their joint density under a local deviation from these values depends only on the $k$ of the largest eigenvalues $\lambda_{1}, \ldots, \lambda_{k}$. Furthermore, we show that $\lambda_{1}, \ldots, \lambda_{k}$ are asymptotically jointly normal, and the statistical experiment of observing all the eigenvalues of the F matrix converges in the Le Cam sense to a Gaussian shift experiment that depends on the asymptotic means and variances of $\lambda_{1}, \ldots, \lambda_{k}$. In particular, the best statistical inference about sufficiently large spikes in the local asymptotic regime is based on the $k$ of the largest eigenvalues only.


Key words: Spiked F-ratio, Local Asymptotic Normality, multivariate F distribution, phase transition, super-critical regime, asymptotic normality of eigenvalues, limits of statistical experiments.

1. Introduction. The roots of the equation

$$
\begin{equation*}
\operatorname{det}(\mathbf{H}-\lambda \mathbf{E})=0, \tag{1}
\end{equation*}
$$

or equivalently the eigenvalues of the $\mathbf{F}$-ratio $\mathbf{E}^{-1} \mathbf{H}$, where matrices $\mathbf{H}$ and $\mathbf{E}$ are the 'hypothesis' and 'error' sums of squares, are fundamental for the multivariate statistics. They form the basis for many invariant tests, including the classical tests of the equality of two covariance matrices and of the linear hypotheses in the multivariate regression. In this paper, we study the behavior of these roots when the F-ratio matrix is high dimensional as is often the case in the contemporaneous statistical applications.

We assume that under the null, both $\mathbf{H}$ and $\mathbf{E}$ are central Wisharts, whereas under the alternative, the 'hypothesis' sum of squares matrix $\mathbf{H}$ contains a low-rank structure. This structure is revealed either in a lowrank difference between the covariance parameters of $\mathbf{H}$ and $\mathbf{E}$, or in a lowrank non-centrality in $\mathbf{H}$. The former corresponds to testing the equality
of two covariance matrices, whereas the latter corresponds to testing linear hypotheses in multivariate regression.

We call the eigenvalues of the low-rank difference between the parameters of $\mathbf{H}$ and $\mathbf{E}$ the spikes. Spiked models have attracted much recent research attention. They were introduced in Johnstone (2001) as a useful abstraction capturing the fact that the high-dimensional sample variation often concentrates along a small number of distinct directions.
In the case of testing the equality of two covariance matrices, these directions may correspond to a few signals that are present only in one of the two samples. We will refer to this as the signal detection case (SigD). In the regression context (REG), an example of a low-rank alternative would be a one-way MANOVA with unequal group means that belong to the same low dimensional hyperplane. Another regression example is the structural break in the number of factors in mean, with a small number of additional factors potentially born by a break event which splits the sample.

The focus of this paper is on the F-ratios of high dimensionality $p$. We consider the asymptotic regime where $p$ goes to infinity proportionally to the 'sample sizes' represented by the 'hypothesis' (matrix H) and 'error' (matrix E) degrees of freedom (d.f.). Our main results can be summarized as follows.

First, we establish a phase transition threshold such that if the spikes are below it, or sub-critical, then any finite number of the largest eigenvalues of the $\mathbf{F}$-ratio almost surely (a.s.) converge to the upper boundary of the support of the limiting spectral distribution of $\mathbf{E}^{-1} \mathbf{H}$, derived by Wachter (1980). In contrast, when $m$ of the spikes are super-critical, the $m$ of the largest eigenvalues of the F-ratio a.s. converge to locations strictly above the upper boundary of the Wachter distribution. The threshold turns out to be the same for SigD and REG cases.

Second, we prove the joint asymptotic normality of the $m$ of the largest eigenvalues of the F-ratio that correspond to the super-critical spikes. We derive explicit formulas for the asymptotic means and variances. In both SigD and REG cases, the asymptotic variance is highly sensitive to the ratio of the dimensionality to the 'error d.f.'. Even small non-zero values of this ratio lead to very substantial variance increases.

Third, and most important, we establish quadratic asymptotic approximations to the likelihood ratios corresponding to local alternatives for $k$ super-critical spikes. We find that the approximations depend only on the $k$ largest eigenvalues of the F-ratio, and that the statistical experiment of observing all the eigenvalues is Locally Asymptotically Normal (LAN). The limiting experiment is a simple $k$-dimensional Gaussian shift. This result im-
plies that the asymptotically optimal inference on the $k$ super-critical spikes can be based exclusively on the $k$ largest eigenvalues of the $\mathbf{F}$-ratio.

We conduct a small-scale Monte Carlo experiment to assess the quality of the LAN confidence sets for super-critical spikes. The experiment shows that the coverage rate of the sets is very close to the nominal one. Moreover, surprisingly, the coverage rate remains good even for low dimensional data.
Iain: I thought it would be nice to add the Monte Carlo results from your second to the last Lumini slide. I have added some experiments and removed reference on bootstrap intervals because I thought that the consensus is that they do not work.

The previous literature on the eigenvalues of $\mathbf{F}$-ratios is vast and old. The finite sample null distribution of the eigenvalues was independently derived by Fisher, Girshick, Hsu, Roy, and Mood in 1939 (see Wilks (1962) for citations). The non-null distributions were classified in James (1964). There have been many subsequent finite sample research papers. These papers are typically motivated by the fact that the power of various tests in MANOVA context depend on the population non-centrality of an F-ratio. To choose between the available tests, one can use the eigenvalues of the $\mathbf{F}$-ratio constructed from a preliminary sample to estimate the non-centrality (see Leung and Muirhead (1987)). As another motivation, Sheena et al (2004) cite the need for estimating the non-centrality of an $\mathbf{F}$-ratio in constructing modified model selection criteria.

In the context of high dimensional data, much recent research focuses on the eigenvalues of sample covariance matrices, which can be viewed as degenerate $\mathbf{F}$-ratios with $\mathbf{E}=I_{p}$. Baik et al (2005) derive the asymptotic distributions of a few of the largest eigenvalues of complex Wisharts. Paul (2007) establishes the asymptotic normality of the fluctuations of a few of the largest eigenvalues of real Wisharts in the super-critical case. Féral and Péché (2009), Benaych-Georges et al (2011) and Bao et al (2014) show that the fluctuations in the sub-critical real case have the Tracy-Widom distribution, while Mo (2012) and Bloemendal and Viràg (2011, 2013) establish the asymptotic distribution of a different type in the critical regime.

In a setting of two independent and not necessarily normal samples with different covariances (SigD case), the phase transition phenomenon has been studied in Nadakuditi and Silverstein (2010). They obtain a formula for the threshold and establish the a.s. limits of the largest eigenvalues corresponding to the super-critical spikes. The asymptotic distribution of the eigenvalues is described in their paper as an open problem. Our paper solves this problem for the case of two normal samples, including the REG case, which was not covered by Nadakuditi and Silverstein (2010).

We expect that our asymptotic normality results can be extended to the F-ratios constructed from non-normal samples. In the one-sample case, an extension of Paul's (2007) asymptotic normality results has been done in Bai and Yao (2008). For non-degenerate F-ratios, our asymptotic normality result for SigD case has been recently extended by Wang and Yao (2015). Iain: Bai and Yao's theorem is incorrect, Wang and Yao's proofs are incomplete. How should we handle this?

The focus of this paper on normal data is dictated by our main goal: establishing the LAN property for the eigenvalues of the $\mathbf{F}$-ratio. To reach this goal, we derive an asymptotic approximation to a log likelihood process by representing it in the form of a multiple contour integral, and applying the Laplace approximation method. The explicit form of the joint distribution of the eigenvalues of $\mathbf{E}^{-1} \mathbf{H}$ is known only in the normal case, and we need such an explicit form for our analysis.

A decision-theoretic approach to the finite sample estimation of the eigenvalues of the "ratio" of the population covariances of $\mathbf{H}$ and $\mathbf{E}$, or the eigenvalues of the non-centrality parameter of $\mathbf{H}$ was taken in many previous studies (see Sheena et al (2004), Bilodeau and Srivastava (1992), and references therein). In one of the first such studies, Muirhead and Verathaworn (1985) explain that the ideal decision-theoretic approach that directly analyzes expected loss with respect to the joint distribution of the eigenvalues of $\mathbf{E}^{-1} \mathbf{H}$ "does not seem feasible due primarily to the complexity of the distribution of the ordered latent roots..." Instead, they focus on deriving an optimal estimator from a particular class.

The proportional asymptotics used in this paper preserves a salient feature of the finite sample, by making the dimensionality of the data non-negligible relative to the sample size. From this perspective, our LAN result can be viewed as an asymptotic implementation of the ideal decision-theoretic approach to the finite sample estimation. We overcome the complexity of the joint distribution of the eigenvalues by using a tractable multiple contour integral representation of the log likelihood process, which follows from the multiple contour integral representation of hypergeometric functions of two matrix arguments, established in Onatski (2013), Dharmawansa and Johnstone (2014), and Passemier et al (2014).

The LAN result of this paper stays in sharp contrast to the asymptotic behavior of the likelihood ratio in the sub-critical regime. In a separate paper, we show that the statistical experiment of observing the eigenvalues of an F-ratio with a single sub-critical spike is not LAN. The corresponding likelihood ratio depends only on a smooth functional of the empirical distribution of all the eigenvalues of $\mathbf{E}^{-1} \mathbf{H}$, so that asymptotically optimal
inference about the spike may ignore information contained in the largest eigenvalue. This is totally different from what happens in the super-critical regime, as our LAN result implies that the asymptotically optimal inference about super-critical spikes can be based on the corresponding largest eigenvalues only.

The rest of the paper is structured as follows. In the next section, we describe our setting. In Section 3, we explore the phase transition and derive the a.s. limits of the super-critical eigenvalues. In Section 4, we establish the asymptotic normality of the super-critical eigenvalues. In Section 5, we derive an asymptotic approximation to the joint distribution of all the eigenvalues of $\mathbf{E}^{-1} \mathbf{H}$ for the case of $k$ super-critical spikes. In Section 6 , we show that the likelihood ratio in the local parameter space is asymptotically equivalent to a linear combination of $k$ of the largest eigenvalues, and establish the LAN property. Section 7 concludes.
2. Setup. Suppose that

$$
\left(n_{1}+k\right) \mathbf{H} \sim W_{p}\left(n_{1}+k, \Sigma_{1}, \Omega_{1}\right) \quad \text { and } \quad n_{2} \mathbf{E} \sim W_{p}\left(n_{2}, \Sigma_{2}\right)
$$

are independent non-central and central Wishart matrices respectively. For the non-centrality parameter $\Omega_{1}$, we use a symmetric version of the definition in Muirhead (1982, p. 442). That is, if $Z$ is an $n \times p$ matrix distributed as $N\left(M, I_{n} \otimes \Sigma\right)$, then $Z^{\prime} Z \sim W_{p}(n, \Sigma, \Omega)$ with the non-centrality parameter $\Omega=\Sigma^{-1 / 2} M^{\prime} M \Sigma^{-1 / 2}$. We are interested in the eigenvalues $\lambda_{p 1} \geq \ldots \geq \lambda_{p p}$ of $\mathbf{F} \equiv \mathbf{E}^{-1} \mathbf{H}$.

In what follows, we will assume that $\Sigma_{2}=I_{p}$. This assumption is without loss of generality because the eigenvalues of $\mathbf{F}$ do not change under the transformation $\mathbf{H} \mapsto \Sigma_{2}^{-1 / 2} \mathbf{H} \Sigma_{2}^{-1 / 2}, \mathbf{E} \mapsto \Sigma_{2}^{-1 / 2} \mathbf{E} \Sigma_{2}^{-1 / 2}$. We will consider two different settings for the parameters $\Sigma_{1}$ and $\Omega_{1}$.

1. Spiked covariance (SigD): $\Sigma_{1}=I_{p}+L s L^{\prime}$ and $\Omega_{1}=0$, where $s=\operatorname{diag}\left\{s_{1}, \ldots, s_{k}\right\}$ with $s_{1}>\ldots>s_{k}>0$ is the diagonal matrix of the covariance spikes, and $L$ is a $p \times k$ matrix with orthogonal columns, which consists of nuisance parameters.
2. Spiked non-centrality (REG): $\Sigma_{1}=I_{p}$ and $\Omega_{1}=n_{H} L s L^{\prime}$ with $n_{H}=n_{1}+k$, where $L$ and $s$ are as defined above, but the diagonal elements of $s$ are interpreted as non-centrality spikes.
Is is convenient to think of $\mathbf{H}$ as the sample covariance matrix $\mathbf{X X}^{\prime} / n_{H}$ of a sample $\mathbf{X}$ having the factor structure

$$
\begin{equation*}
\mathbf{X}=L F^{\prime}+\varepsilon \tag{2}
\end{equation*}
$$

with $L, F$, and $\varepsilon$ playing the roles of the normalized factor loadings, factors, and idiosyncratic terms, respectively. Matrices $F$ and $\varepsilon$ are mutually independent, and independent from $\mathbf{E}$. The entries of $\varepsilon$ are i.i.d. standard normals, and the distribution of $F$ depends on the setting. For SigD, $F \sim N\left(0, I_{n_{H}} \otimes s\right)$, whereas for REG, $F$ is a deterministic matrix such that $F^{\prime} F / n_{H}=s$. With this interpretation, SigD and REG describe, respectively, distributions of $\mathbf{H}$ which are unconditional and conditional on the factors. In both cases the spike parameters $s_{j}$ measure the $j$-th factor's variability or 'strength'.

Let us introduce a convenient representation for the eigenvalues of $\mathbf{F}$. First, note that these eigenvalues are invariant with respect to the simultaneous transformations

$$
\begin{equation*}
\mathbf{X} \mapsto U \mathbf{X} V \equiv X \quad \text { and } \quad \mathbf{E} \mapsto U \mathbf{E} U^{\prime} \equiv E \tag{3}
\end{equation*}
$$

where $U$ is a random matrix uniformly distributed over the orthogonal group $\mathcal{O}(p)$, and $V \in \mathcal{O}\left(n_{H}\right)$ is such that the submatrix of its first $k$ columns equals $F\left(F^{\prime} F\right)^{-1 / 2}$.

Note that $U L$ can be represented as

$$
U L=v\left(v^{\prime} v\right)^{-1 / 2} \equiv v W_{v}^{-1 / 2}
$$

where $v$ is a $p \times k$ matrix with i.i.d. standard normal entries, and $W_{v} \sim$ $W_{k}\left(p, I_{k}\right)$. Furthermore,

$$
X=\left[v W_{v}^{-1 / 2} s^{1 / 2} W_{F}^{1 / 2}, 0\right]+\epsilon,
$$

where $v, \epsilon$, and $W_{F}$ are mutually independent, the entries of $\epsilon$ are i.i.d. standard normals, and the distribution of $W_{F}$ depends on the setting. For SigD, $W_{F} \sim W_{k}\left(n_{H}, I_{k}\right)$, whereas for REG, $W_{F}=n_{H} I_{k}$.

Let us denote the submatrix of the first $k$ columns of $\epsilon$ as $u$. Then

$$
\begin{equation*}
X X^{\prime}=\xi \xi^{\prime}+n_{1} H \tag{4}
\end{equation*}
$$

where $n_{1} H \sim W_{p}\left(n_{1}, I_{p}\right), H$ and $\xi \xi^{\prime}$ are mutually independent, and independent from $E$, and

$$
\begin{equation*}
\xi=v W_{v}^{-1 / 2} s^{1 / 2} W_{F}^{1 / 2}+u . \tag{5}
\end{equation*}
$$

Using (3) and (4), we obtain the following convenient representation for the eigenvalues $\lambda_{p 1} \geq \ldots \geq \lambda_{p p}$ of $\mathbf{F}$. Let $\hat{\lambda}_{p 1} \geq \ldots \geq \hat{\lambda}_{p p}$ be the roots of the equation

$$
\begin{equation*}
\operatorname{det}\left(\xi \xi^{\prime} / n_{1}+H-x E\right)=0 \tag{6}
\end{equation*}
$$

Then

$$
\begin{equation*}
\lambda_{p j}=n_{1} \hat{\lambda}_{p j} /\left(n_{1}+k\right) . \tag{7}
\end{equation*}
$$

This representation is convenient because the roots of (6) can be viewed and analyzed as perturbations of the roots of equation $\operatorname{det}(H-x E)=0$ caused by adding the low-rank matrix $\xi \xi^{\prime} / n_{1}$ to $H$. Here $\xi, H$, and $E$ are independent and

$$
n_{1} H \sim W_{p}\left(n_{1}, I_{p}\right), \quad n_{2} E \sim W_{p}\left(n_{2}, I_{p}\right)
$$

If $x \in \mathbb{R}$ is such that $H-x E$ is invertible, then

$$
\left(\xi \xi^{\prime} / n_{1}+H-x E\right)^{-1}=Q-Q \xi\left(I_{k}+\xi^{\prime} Q \xi / n_{1}\right)^{-1} \xi^{\prime} Q / n_{1}
$$

where $Q \equiv(H-x E)^{-1}$. Therefore, if $x$ is a root of the equation

$$
\begin{equation*}
\operatorname{det}\left(I_{k}+\xi^{\prime}(H-x E)^{-1} \xi / n_{1}\right)=0 \tag{8}
\end{equation*}
$$

then it also solves (6), and hence, behavior of the roots of (6) can be inferred from that of the random matrix-valued function

$$
\begin{equation*}
M(x)=\xi^{\prime}(H-x E)^{-1} \xi / n_{1} . \tag{9}
\end{equation*}
$$

This is the main idea of the analysis in the next section.
3. Phase transition and almost sure limits. We will consider the proportional asymptotic regime where $n_{1}, n_{2}$, and $p$ diverge to infinity so that

$$
c_{1} \equiv p / n_{1} \rightarrow \gamma_{1} \text { and } c_{2} \equiv p / n_{2} \rightarrow \gamma_{2} \text { with } \gamma_{j} \in(0,1)
$$

Let $n=\left(n_{1}, n_{2}\right)$ and $\gamma=\left(\gamma_{1}, \gamma_{2}\right)$. We will abbreviate the above asymptotics as $p, n \rightarrow \gamma$.

As follows from Wachter's (1980) work (see also Yin et al. (1983) and Silverstein (1985)), as $p, n \rightarrow_{\gamma} \infty$, the empirical distribution of the eigenvalues of $E^{-1} H$ converges in probability to the distribution with density

$$
\begin{equation*}
\frac{1-\gamma_{2}}{2 \pi} \frac{\sqrt{\left(\beta_{+}-\lambda\right)\left(\lambda-\beta_{-}\right)}}{\lambda\left(\gamma_{1}+\gamma_{2} \lambda\right)} \mathbf{1}\left\{\beta_{-} \leq \lambda \leq \beta_{+}\right\} . \tag{10}
\end{equation*}
$$

The upper and the lower boundaries of the support of this density are

$$
\beta_{ \pm}=\left(\frac{1 \pm \rho}{1-\gamma_{2}}\right)^{2}, \text { where } \rho=\sqrt{\gamma_{1}+\gamma_{2}-\gamma_{1} \gamma_{2}} .
$$

The results of Silverstein and Bai (1995) and Silverstein (1995) show that the empirical distribution converges not only in probability, but also a.s. Furthermore, a simple extension of Theorem 1.1 of Bai and Silverstein (1998) that covers random $E$ (Lemma SM1 in the Supplementary Material) implies that the largest eigenvalue of $E^{-1} H$ a.s. converges to $\beta_{+}$.

The latter convergence, together with (7) and Weyl's inequalities for the eigenvalues of a sum of two Hermitian matrices (see Theorem 4.3.7 in Horn and Johnson (1985)), imply that the $k+1$-th largest eigenvalue of $\mathbf{F}, \lambda_{p, k+1}$, a.s. converges to $\beta_{+}$. Those of the $k$ largest eigenvalues that remain separated from $\beta_{+}$as $p, n \rightarrow_{\gamma} \infty$, must correspond to solutions of (8). Below, we study these solutions in detail. The Supplementary Material (SM) contains proofs of the following three auxiliary lemmas.

Lemma 1. For any $x>\beta_{+}$, as $p, n \rightarrow_{\gamma} \infty$,

$$
\begin{align*}
& \operatorname{tr}(H-x E)^{-1} / p \xrightarrow{\text { a.s. }} m(0 ; x) \text { and }  \tag{11}\\
& \frac{\mathrm{d}}{\mathrm{~d} x} \operatorname{tr}(H-x E)^{-1} / p \xrightarrow{\text { a.s. }} \frac{\mathrm{d}}{\mathrm{~d} x} m(0 ; x), \tag{12}
\end{align*}
$$

where $m(0 ; x)=\lim _{z \rightarrow 0} m(z ; x)$, and $m \equiv m(z ; x) \in \mathbb{C}^{+}$is an analytic function of $z \in \mathbb{C}^{+}$that satisfies equation

$$
\begin{equation*}
z-\frac{1}{1+\gamma_{1} m}=-\frac{1}{m}-\frac{x}{1-\gamma_{2} x m} \tag{13}
\end{equation*}
$$

Lemma 2. For any $x>\beta_{+}$, as $p, n \rightarrow_{\gamma} \infty$,

$$
\begin{aligned}
& \left\|M(x)-\left(s+\gamma_{1} I_{k}\right) \operatorname{tr}(H-x E)^{-1} / p\right\| \xrightarrow{\text { a.s. }} 0 \text { and } \\
& \left\|\frac{\mathrm{d}}{\mathrm{~d} x} M(x)-\left(s+\gamma_{1} I_{k}\right) \frac{\mathrm{d}}{\mathrm{~d} x} \operatorname{tr}(H-x E)^{-1} / p\right\| \xrightarrow{\text { a.s. }} 0,
\end{aligned}
$$

where $\|\cdot\|$ denotes the spectral norm.
In the next lemma, and throughout this paper, the statement "for sufficiently large $p$ " abbreviates "for sufficiently large $p$ and $n$ along the sequence $p, n \rightarrow_{\gamma} \infty$.

Lemma 3. (i) For any $\varepsilon>0$, the eigenvalues of $M(x)$ are strictly increasing functions of $x \in\left(\beta_{+}+\varepsilon, \infty\right)$ for sufficiently large $p$, a.s.;
(ii) $m(0 ; x)$ is a strictly increasing, continuous function of $x \in\left(\beta_{+}, \infty\right)$;
(iii) $\lim _{x \rightarrow \infty} m(0 ; x)=0$, and $\lim _{x \downarrow \beta_{+}} m(0 ; x)\left(s_{i}+\gamma_{1}\right)<-1$ if and only if $s_{i}>\bar{s}$, where

$$
\bar{s}=\left(\gamma_{2}+\rho\right) /\left(1-\gamma_{2}\right)
$$

Let $\hat{\lambda}_{p 1} \geq \ldots \geq \hat{\lambda}_{p k}$ be the $k$ largest solutions of equation (8). By Lemmas $1-3$, if

$$
\begin{equation*}
s_{1}>\ldots>s_{q}>\bar{s}>s_{q+1}>\ldots>s_{k} \tag{14}
\end{equation*}
$$

then $\hat{\lambda}_{p i} \xrightarrow{\text { a.s. }} x_{i}$, where $x_{i}, i=1, \ldots, q$, are such that

$$
\begin{equation*}
1+\left(s_{i}+\gamma_{1}\right) m\left(0 ; x_{i}\right)=0 \tag{15}
\end{equation*}
$$

and $m\left(0 ; x_{i}\right)$ satisfies (13) with $x$ replaced by $x_{i}$. In particular,

$$
\begin{equation*}
\frac{1}{1+\gamma_{1} m\left(0 ; x_{i}\right)}-\frac{1}{m\left(0 ; x_{i}\right)}-\frac{x_{i}}{1-\gamma_{2} x_{i} m\left(0 ; x_{i}\right)}=0 \tag{16}
\end{equation*}
$$

Combining (15) and (16), we obtain

$$
\frac{1}{s_{i}}+1-\frac{x_{i}}{s_{i}+\gamma_{1}+\gamma_{2} x_{i}}=0
$$

which implies that

$$
\begin{equation*}
x_{i}=\frac{\left(s_{i}+\gamma_{1}\right)\left(s_{i}+1\right)}{s_{i}-\gamma_{2}\left(s_{i}+1\right)} \tag{17}
\end{equation*}
$$

By (7), $n_{1} \hat{\lambda}_{p i} /\left(n_{1}+k\right), i=1, \ldots, q$, must be the $q$ largest eigenvalues of $\mathbf{F}$, and thus, $x_{i}, i=1, \ldots, q$, describe their a.s. limits. Since there are only $q$ roots of (8) that are asymptotically separated from $\beta_{+}$and are located above $\beta_{+}$, the other $k-q$ of the largest eigenvalues of $\mathbf{F}$ must a.s. converge to $\beta_{+}$. To summarize, the following theorem holds.

TheOrem 4. Suppose that $s_{1}>\ldots>s_{q}>\bar{s}>s_{q+1}>\ldots>s_{k}$, and let $\lambda_{p i}$ be the $i$-th largest eigenvalue of $\mathbf{F}$. Then for $i \leq q$,

$$
\lambda_{p i} \xrightarrow{\text { a.s. }} \frac{\left(s_{i}+\gamma_{1}\right)\left(s_{i}+1\right)}{s_{i}-\gamma_{2}\left(s_{i}+1\right)}
$$

as $p, n \rightarrow{ }_{\gamma} \infty$. For $q<i \leq k, \lambda_{p i} \xrightarrow{\text { a.s. }} \beta_{+}$.
As follows from Theorem $4, \bar{s}=\left(\gamma_{2}+\rho\right) /\left(1-\gamma_{2}\right)$ is the phase transition threshold for the eigenvalues of the spiked $\mathbf{F}$-ratio. The value of this threshold diverges to infinity when $\gamma_{2} \rightarrow 1$. Note that when $\gamma_{2}$ is close to one, the smallest eigenvalue of $E$ is close to zero, which makes $E^{-1}$ a particularly bad estimator of the inverse of the population covariance $\Sigma_{2}^{-1}$. When $\gamma_{2} \rightarrow 0$, the threshold converges to $\sqrt{\gamma_{1}}$, which is the phase transition threshold for
the eigenvalues of one spiked Wishart matrix. In such a case, $x_{i}$ converges to $\left(s_{i}+\gamma_{1}\right)\left(s_{i}+1\right) / s_{i}$, which is the a.s. limit of the $i$-th largest eigenvalue of the spiked Wishart when the $i$-th spike $s_{i}$ is above $\sqrt{\gamma_{1}}$.

When both $\gamma_{1}$ and $\gamma_{2}$ converge to zero, $x_{i}$ converges to $s_{i}+1$, which is the population analogue of $\lambda_{p i}$. For positive $\gamma_{1}$ and $\gamma_{2}, \lambda_{p i}$ is an upward biased estimator of $s_{i}+1$. The relative bias $\lambda_{p i} /\left(s_{i}+1\right)$ converges to $1 /\left(1-\gamma_{2}\right)$ when the spike $s_{i}$ diverges to infinity. The sizes of the relative and absolute biases are very senistive to the value of $\gamma_{2}$. They quickly increase when $\gamma_{2}$ rises above zero. Such a behavior is illustrated in Figure 1.


FIg 1. The a.s. limits of $\lambda_{p i}$ as functions of $s_{i}$ for $s_{i}>\bar{s}$. The top line corresponds to $\gamma_{1}=1 / 2, \gamma_{2}=1 / 10$. The dashed line corresponds to the same $\gamma_{1}$, but $\gamma_{2}=0$. The straight line is the population analogue of $\lambda_{p i}, s_{i}+1$.
4. Asymptotic normality. In this section, we will assume that (14) holds, so that only $q$ eigenvalues of $\mathbf{F}$ separate from the bulk asymptotically. We would like to study their fluctuations around the corresponding a.s. limits. Theorem 4 shows that the limits $x_{i}$ depend on $\gamma_{1}$ and $\gamma_{2}$. Because of this dependence, the rate of the convergence has to depend on the rates of the convergences $c_{1} \rightarrow \gamma_{1}$ and $c_{2} \rightarrow \gamma_{2}$. However, as will be shown below, the latter rates do not affect the fluctuations of $\lambda_{i}$ around

$$
l_{i}=\frac{\left(s_{i}+c_{1}\right)\left(s_{i}+1\right)}{s_{i}-c_{2}\left(s_{i}+1\right)}
$$

which are obtained from $x_{i}$ by replacing $\gamma_{1}$ and $\gamma_{2}$ by $c_{1} \equiv p / n_{1}$ and $c_{2} \equiv$ $p / n_{2}$ in equation (17).

Similar to $x_{i}$, which are linked to the Stieltjes transform of the limiting spectral distribution of $H-x E$ via (15), $l_{i}$ also can be linked to the limiting Stieltjes transform, albeit under a slightly different asymptotic regime. Precisely, let $m_{p}(z ; x)$ be the Stieltjes transform of the limiting spectral distribution of $H-x E$ as $n_{1}, n_{2}$, and $p$ diverge to infinity so that $p / n_{1}$ and $p / n_{2}$ remain fixed. Then, similarly to (15), we have

$$
\begin{equation*}
1+\left(s_{i}+c_{1}\right) m_{p}\left(0 ; l_{i}\right)=0 . \tag{18}
\end{equation*}
$$

This equation will be useful in our analysis below, where we maintain the assumption that $p / n_{1}$ and $p / n_{2}$ are not necessarily fixed, but converge to $\gamma_{1}$ and $\gamma_{2}$, respectively.

Recall that, by (7), $\lambda_{p i}=n_{1} \hat{\lambda}_{p i} /\left(n_{1}+k\right)$, where $\hat{\lambda}_{p i}, i=1, \ldots, q$, satisfy (8). Clearly, the asymptotic distributions of $\sqrt{p}\left(\lambda_{p i}-l_{i}\right)$ and $\sqrt{p}\left(\hat{\lambda}_{p i}-l_{i}\right)$, $i=1, \ldots, q$, coincide. Therefore, below we will study the asymptotic behavior of the latter. By the standard Taylor expansion argument,

$$
\begin{equation*}
\sqrt{p}\left(\hat{\lambda}_{p i}-l_{i}\right)=-\frac{\sqrt{p} \operatorname{det} \mathcal{M}\left(l_{i}\right)}{\frac{\mathrm{d}}{\mathrm{~d} x} \operatorname{det} \mathcal{M}\left(l_{i}\right)+\frac{1}{2}\left(\hat{\lambda}_{p i}-l_{i}\right) \frac{\mathrm{d}^{2}}{\mathrm{~d} x^{2}} \operatorname{det} \mathcal{M}\left(\tilde{l}_{i}\right)}, \tag{19}
\end{equation*}
$$

$i=1, \ldots, q$, where $\mathcal{M}(x)=I_{k}+M(x)$, and $\tilde{l}_{i} \in\left[l_{i}, \hat{\lambda}_{p i}\right]$.
We have (see, for example, Magnus and Neudecker (1999) pp. 149-150)

$$
\begin{aligned}
\frac{\mathrm{d}}{\mathrm{~d} x} \operatorname{det} \mathcal{M}\left(l_{i}\right) & =\operatorname{det} \mathcal{M}\left(l_{i}\right) \operatorname{tr} J\left(l_{i}\right), \text { and } \\
\frac{\mathrm{d}^{2}}{\mathrm{~d} x^{2}} \operatorname{det} \mathcal{M}\left(l_{i}\right) & =\operatorname{det} \mathcal{M}\left(l_{i}\right)\left\{\operatorname{tr} R\left(l_{i}\right)+\left(\operatorname{tr} J\left(l_{i}\right)\right)^{2}-\operatorname{tr} J^{2}\left(l_{i}\right)\right\}
\end{aligned}
$$

where

$$
J(x)=\mathcal{M}(x)^{-1} \frac{\mathrm{~d}}{\mathrm{~d} x} M(x), \text { and } R(x)=\mathcal{M}(x)^{-1} \frac{\mathrm{~d}^{2}}{\mathrm{~d} x^{2}} M(x) .
$$

Since the event

$$
\operatorname{det} \mathcal{M}\left(l_{i}\right)=0 \text { or } 1+M_{i i}\left(l_{i}\right)=0 \text { for some } i=1, \ldots, q
$$

happens with probability zero, we can simultaneously multiply the numerator and denominator of (19) by $\left(1+M_{i i}\left(l_{i}\right)\right) / \operatorname{det} \mathcal{M}\left(l_{i}\right)$ to obtain

$$
\begin{equation*}
\sqrt{p}\left(\hat{\lambda}_{p i}-l_{i}\right)=-\frac{\sqrt{p}\left(1+M_{i i}\left(l_{i}\right)\right)}{\sigma\left(l_{i}\right)+\frac{1}{2}\left(\hat{\lambda}_{p i}-l_{i}\right) \delta\left(l_{i}\right)} \tag{20}
\end{equation*}
$$

where

$$
\begin{aligned}
\sigma\left(l_{i}\right) & =\left(1+M_{i i}\left(l_{i}\right)\right) \operatorname{tr} J\left(l_{i}\right), \text { and } \\
\delta\left(l_{i}\right) & =\left(1+M_{i i}\left(l_{i}\right)\right)\left\{\operatorname{tr} R\left(l_{i}\right)+\left(\operatorname{tr} J\left(l_{i}\right)\right)^{2}-\operatorname{tr} J^{2}\left(l_{i}\right)\right\} .
\end{aligned}
$$

A proof of the following lemma is given in the SM.
Lemma 5. For any $i=1, \ldots$, $q$, we have: $(i) ~ \sigma\left(l_{i}\right) \xrightarrow{\mathrm{P}}\left(s_{i}+\gamma_{1}\right) \frac{\mathrm{d}}{\mathrm{d} x} m\left(0 ; x_{i}\right)$; (ii) $\delta\left(l_{i}\right)=O(1)$ a.s.

Equation (20), Lemma 5, and the Slutsky theorem imply that, for the purpose of establishing convergence in distribution of $\sqrt{p}\left(\hat{\lambda}_{p i}-l_{i}\right), i=$ $1, \ldots, q$, we may focus on the numerator of (20)

$$
Z_{i i}\left(l_{i}\right) \equiv \sqrt{p}\left(1+M_{i i}\left(l_{i}\right)\right)=\sqrt{p}\left[M_{i i}\left(l_{i}\right)-\left(s_{i}+c_{1}\right) m_{p}\left(0 ; l_{i}\right)\right],
$$

where the last equality follows from (18).
The random variable $Z_{i i}$ is the entry of the matrix

$$
Z\left(l_{i}\right)=\sqrt{p}\left[M\left(l_{i}\right)-\left(s+c_{1} I_{k}\right) m_{p}\left(0 ; l_{i}\right)\right]
$$

that belongs to the $i$-th row and the $i$-th column. Let us now introduce new notations. Let

$$
\begin{aligned}
D & =\left(W_{F} / n_{1}\right)^{1 / 2} s^{1 / 2}\left(W_{v} / p\right)^{-1 / 2}, \\
G & =\left(H-l_{i} E\right)^{-1} / p, \\
\Delta_{F} & =\sqrt{n_{1}}\left(\left(W_{F} / n_{1}\right)^{1 / 2}-I_{k}\right), \text { and } \\
\Delta_{v} & =\sqrt{p}\left(W_{v} / p-I_{k}\right) .
\end{aligned}
$$

Then, using equations (9) and (5), we obtain the following decomposition.

$$
Z\left(l_{i}\right)=\sum_{v=1}^{7} Z^{(v)},
$$

where

$$
\begin{aligned}
& Z^{(1)}=D \sqrt{p}\left(v^{\prime} G v-I_{k} \operatorname{tr} G\right) D^{\prime}, \\
& Z^{(2)}=(\operatorname{tr} G) D\left(W_{v} / p\right)^{-1 / 2} s^{1 / 2} \sqrt{c_{1}} \Delta_{F}, \\
& Z^{(3)}=\operatorname{tr} G \sqrt{c_{1}} \Delta_{F} s^{1 / 2}\left(W_{v} / p\right)^{-1} s^{1 / 2}, \\
& Z^{(4)}=-(\operatorname{tr} G) s^{1 / 2} \Delta_{v}\left(W_{v} / p\right)^{-1} s^{1 / 2}, \\
& Z^{(5)}=\sqrt{c_{1}} \sqrt{p}\left(D v^{\prime} G u+u^{\prime} G v D^{\prime}\right), \\
& Z^{(6)}=c_{1} \sqrt{p}\left(u^{\prime} G u-I_{k} \operatorname{tr} G\right),
\end{aligned}
$$

and

$$
Z^{(7)}=\left(s+c_{1} I_{k}\right) \sqrt{p}\left(\operatorname{tr} G-m_{p}\left(0 ; l_{i}\right)\right) .
$$

For the last term, $Z^{(7)}$, we have the following lemma.
Lemma 6. $\quad Z(7) \xrightarrow{\text { a.s. }} 0$.
A proof of this lemma is given in the SM. Had $l_{i}$ been negative, $H-l_{i} E$ would have been having the form $Y T Y^{\prime}$ with $Y \sim N\left(0, I_{p} \otimes I_{n_{1}+n_{2}}\right)$ and a positive definite diagonal $T$ with converging spectral distribution. Then Lemma 6 would have been following from the results of Bai and Silverstein (2004). Our proof extends Bai and Silverstein's (2004) arguments to the case of negative $l_{i}$.

Further, the asymptotic behavior of the terms $Z^{(2)}$ and $Z^{(3)}$ differ depending on the setting. Recall that for SigD, $W_{F} \sim W_{k}\left(n_{H}, I_{k}\right)$. Then, since

$$
\Delta_{F}=\sqrt{n_{1}}\left(W_{F} / n_{1}-I_{k}\right) / 2+o_{\mathrm{P}}(1)
$$

a standard CLT together with Lemma 1 imply that

$$
\begin{equation*}
\operatorname{diag}\left(Z^{(2)}+Z^{(3)}\right) \xrightarrow{d} N\left(0,2 \gamma_{1} m^{2}\left(0 ; x_{i}\right) s^{2}\right) \tag{21}
\end{equation*}
$$

The latter limit is independent from the limits of $Z^{(j)}, j \neq 2,3$, because $W_{F}$ is independent from $u$ and $v$.

In contrast, for REG, we have $W_{F}=n_{H} I_{k}$, and $\Delta_{F}=o(1)$. Therefore,

$$
\begin{equation*}
\operatorname{diag}\left(Z^{(2)}+Z^{(3)}\right) \xrightarrow{\mathrm{P}} 0 \tag{22}
\end{equation*}
$$

Let us now establish the convergence of the remaining components $Z^{(j)}$. Let $b_{i}$ and $B_{i}$ be such that $\left[b_{i}, B_{i}\right]$ includes the support of the limiting spectral distribution, $F\left(\lambda ; x_{i}\right)$, of $H-l_{i} E$. Moreover, let $\left[b_{i}, B_{i}\right]$ be such that none of the eigenvalues $\mu_{1, i} \geq \ldots \geq \mu_{p, i}$ of $H-l_{i} E$ lies outside $\left[b_{i}, B_{i}\right]$ for sufficiently large $p$, a.s. Further, let $g_{j}$ with $j=1, \ldots, J$, where $J$ is an arbitrary positive integer, be functions which are continuous on $\left[b_{i}, B_{i}\right]$ and let $\zeta$ denote a $p \times q$ matrix with i.i.d. $N(0,1)$ entries, independent from $H$ and $E$. Finally, let

$$
\Theta=\{(j, s, t): j=1, \ldots, J ; 1 \leq s \leq t \leq q\} .
$$

The following lemma can be viewed as a special case of Theorem 5.2 in Capitaine et al (2009) or Theorem 7.2 of Bai and Yao (2008), modified to fit the needs of this paper. For readers' convenience, its proof is given in the SM.

Lemma 7. The joint distribution of random variables

$$
\left\{\frac{1}{\sqrt{p}} \sum_{r=1}^{p} g_{j}\left(\mu_{r, i}\right)\left(\zeta_{r s} \zeta_{r t}-\delta_{s t}\right),(j, s, t) \in \Theta\right\}
$$

weakly converges to a multivariate normal. The covariance between components $(j, s, t)$ and $\left(j_{1}, s_{1}, t_{1}\right)$ of the limiting distribution is equal to 0 when $(s, t) \neq\left(s_{1}, t_{1}\right)$, and to $\left(1+\delta_{s t}\right) \int g_{j}(\lambda) g_{j_{1}}(\lambda) \mathrm{d} F\left(\lambda ; x_{i}\right)$ when $(s, t)=\left(s_{1}, t_{1}\right)$.

Note that all entries of $Z^{(j)}, j=1,4,5,6$ are linear combinations of the terms having the form considered in Lemma 7, with weights converging in probability to finite constants. Take for example $Z^{(1)}$. Its entries are linear combinations of the entries of

$$
\frac{1}{\sqrt{p}} v^{\prime}\left(H-l_{i} E\right)^{-1} v-I_{k} \frac{1}{\sqrt{p}} \operatorname{tr}\left(H-l_{i} E\right)^{-1}
$$

which, in turn, can be represented in the form $\frac{1}{\sqrt{p}} \sum_{r=1}^{p}\left(\mu_{r, i}\right)^{-1}\left(\zeta_{r s} \zeta_{r t}-\delta_{s t}\right)$. The matrix $\zeta$ is obtained by multiplying $[u, v]$ from the left by the eigenvector matrix of $H-l_{i} E$.
Lemma 7 implies that vector $\left(Z_{i i}^{(1)}, Z_{i i}^{(4)}, Z_{i i}^{(5)}, Z_{i i}^{(6)}\right)$ converges in distribution to a four-dimensional normal vector with zero mean and the following covariance matrix

$$
\left(\begin{array}{cccc}
2 s_{i}^{2} m^{\prime}\left(0 ; x_{i}\right) & -2 s_{i}^{2} m^{2}\left(0 ; x_{i}\right) & 0 & 0 \\
-2 s_{i}^{2} m^{2}\left(0 ; x_{i}\right) & 2 s_{i}^{2} m^{2}\left(0 ; x_{i}\right) & 0 & 0 \\
0 & 0 & 4 \gamma_{1} s_{i} m^{\prime}\left(0 ; x_{i}\right) & 0 \\
0 & 0 & 0 & 2 \gamma_{1}^{2} m^{\prime}\left(0 ; x_{i}\right)
\end{array}\right) .
$$

Combining this result with Lemma 6, and convergencies (21) and (22), we obtain, for SigD,

$$
\begin{equation*}
Z_{i i}\left(l_{i}\right) \xrightarrow{d} N\left(0,2\left(s_{i}+\gamma_{1}\right)^{2} m^{\prime}\left(0 ; x_{i}\right)-2 s_{i}^{2}\left(1-\gamma_{1}\right) m^{2}\left(0 ; x_{i}\right)\right), \tag{23}
\end{equation*}
$$

and, for REG,

$$
\begin{equation*}
Z_{i i}\left(l_{i}\right) \xrightarrow{d} N\left(0,2\left(s_{i}+\gamma_{1}\right)^{2} m^{\prime}\left(0 ; x_{i}\right)-2 s_{i}^{2} m^{2}\left(0 ; x_{i}\right)\right) . \tag{24}
\end{equation*}
$$

To establish the joint convergence of $Z_{i i}\left(l_{i}\right), i=1, \ldots, q$, we need another lemma. For each $i=1, \ldots, q$, let $g_{j, i}$, with $j=1, \ldots, J$, be functions continuous on $\left[b_{i}, B_{i}\right]$.

Lemma 8. For any set of pairs $\left\{\left(s_{i}, t_{i}\right): i=1, \ldots, q\right\}$ such that $\left(s_{i_{1}}, t_{i_{1}}\right) \neq$ $\left(s_{i_{2}}, t_{i_{2}}\right)$ for any $i_{1} \neq i_{2}$, the joint distribution of random variables

$$
\left\{\frac{1}{\sqrt{p}} \sum_{r=1}^{p} g_{j, i}\left(\mu_{r, i}\right)\left(\zeta_{r s_{i}} \zeta_{r t_{i}}-\delta_{s_{i} t_{i}}\right), i=1, \ldots, q\right\}
$$

weakly converges to a multivariate normal. The covariance between components $i_{1}$ and $i_{2}$ of the limiting distribution is equal to 0 when $i_{1} \neq i_{2}$.

A proof of this lemma is very similar to that of Lemma 7, and we do not report it. Lemma 8 implies that $Z_{i i}\left(l_{i}\right), i=1, \ldots, q$ jointly converge to a $q$-dimensional normal vector with a diagonal covariance matrix. This result, together with equations (7, 20), Lemma 5, and convergences (23, 24) establish the following lemma.

Lemma 9. The joint asymptotic distribution of $\sqrt{p}\left(\lambda_{p i}-l_{i}\right), i=1, \ldots, q$ is normal, with diagonal covariance matrix. For SigD, the $i$-th diagonal element of the covariance matrix equals

$$
\begin{equation*}
\frac{2\left(s_{i}+\gamma_{1}\right)^{2} m^{\prime}\left(0 ; x_{i}\right)-2 s_{i}^{2}\left(1-\gamma_{1}\right) m^{2}\left(0 ; x_{i}\right)}{\left(s_{i}+\gamma_{1}\right)^{2}\left(\frac{\mathrm{~d}}{\mathrm{~d} x} m\left(0 ; x_{i}\right)\right)^{2}} \tag{25}
\end{equation*}
$$

For $R E G$, it equals

$$
\begin{equation*}
\frac{2\left(s_{i}+\gamma_{1}\right)^{2} m^{\prime}\left(0 ; x_{i}\right)-2 s_{i}^{2} m^{2}\left(0 ; x_{i}\right)}{\left(s_{i}+\gamma_{1}\right)^{2}\left(\frac{\mathrm{~d}}{\mathrm{~d} x} m\left(0 ; x_{i}\right)\right)^{2}} \tag{26}
\end{equation*}
$$

In the SM , we establish the following explicit expressions for $m^{2}\left(0 ; x_{i}\right)$, $m^{\prime}\left(0 ; x_{i}\right)$, and $\frac{\mathrm{d}}{\mathrm{d} x} m\left(0 ; x_{i}\right)$ :

$$
\begin{equation*}
m^{2}\left(0 ; x_{i}\right)=\left(s_{i}+\gamma_{1}\right)^{-2} \tag{27}
\end{equation*}
$$

$$
\begin{equation*}
m^{\prime}\left(0 ; x_{i}\right)=-\frac{s_{i}^{2}}{\left(s_{i}+\gamma_{1}\right)^{2}\left(\gamma_{1}+\gamma_{2}\left(1+s_{i}\right)^{2}-s_{i}^{2}\right)} \tag{28}
\end{equation*}
$$

$$
\begin{equation*}
\mathrm{d} m\left(0 ; x_{i}\right) / \mathrm{d} x=\frac{-\left(\gamma_{2}\left(1+s_{i}\right)-s_{i}\right)^{2}}{\left(s_{i}+\gamma_{1}\right)^{2}\left(\gamma_{1}+\gamma_{2}\left(1+s_{i}\right)^{2}-s_{i}^{2}\right)} \tag{29}
\end{equation*}
$$

Using (27), (28), and (29) in (25) and (26), we obtain

THEOREM 10. Let $\bar{s}=\left(\gamma_{2}+\rho\right) /\left(1-\gamma_{2}\right), \underline{s}=-\left(\gamma_{1}+\gamma_{2}\right) /\left(\gamma_{2}+\rho\right)$, and $d\left(s_{i}\right)=\left(1-\gamma_{2}\right) s_{i}-\gamma_{2}$. Then, for any $s_{1}>\ldots>s_{q}>\bar{s}$, the joint asymptotic distribution of $\sqrt{p}\left(\lambda_{p i}-l_{i}\right), i=1, \ldots, q$ is normal with diagonal covariance matrix. In particular,

$$
\begin{equation*}
\sqrt{p}\left(\lambda_{p i}-l_{i}\right) \xrightarrow{d} N\left(0, \tau^{2}\left(s_{i}\right)\right), \tag{30}
\end{equation*}
$$

where $\tau^{2}\left(s_{i}\right)=w\left(s_{i}\right) \times t\left(s_{i}\right)$ with

$$
t\left(s_{i}\right)=\left(1-\gamma_{2}\right)\left(s_{i}-\underline{s}\right)\left(s_{i}-\bar{s}\right) / d^{2}\left(s_{i}\right)
$$

and

$$
w\left(s_{i}\right)= \begin{cases}2\left(\rho s_{i}\left(1+s_{i}\right) / d\left(s_{i}\right)\right)^{2} & \text { for SigD }  \tag{31}\\ 2\left(\rho s_{i}\left(1+s_{i}\right) / d\left(s_{i}\right)\right)^{2}-2 \gamma_{1} s_{i}^{2} t\left(s_{i}\right) & \text { for REG }\end{cases}
$$

REMARK 11. The $w\left(s_{i}\right)$ component of the asymptotic variance will play the role of the scaling factor in our LAN result below. The fact that the asymptotic variance is smaller for $R E G$ than for SigD accords with intuition. Indeed, as discussed above, REG corresponds to the analysis conditional on factors $F$, whereas SigD corresponds to the unconditional analysis. The factors' variance adds to the asymptotic variance of $\lambda_{p i}$.

Similarly to the bias discussed in the previous section, the asymptotic variance of $\lambda_{p i}$ is sensitive to the size of $\gamma_{2}$. Figure 2 shows that even a small increase in $\gamma_{2}$ may lead to a large increase in the variance.

As the value of the spike $s_{i}$ approaches the phase transition threshold $\bar{s}$ from above, the asymptotic variance converges to zero. As $s_{i} \rightarrow \infty$, the standard deviation increases linearly in $s_{i}$ so that the coefficient of variation does not approach zero for large spikes. The limit of the squared coefficient of variation equals

$$
\lim _{s_{i} \rightarrow \infty} C V^{2}\left(\lambda_{p i}\right)= \begin{cases}2 \rho^{2} /\left(1-\gamma_{2}\right) & \text { for SigD } \\ 2 \gamma_{2} /\left(1-\gamma_{2}\right) & \text { for REG }\end{cases}
$$

Again, it is sensitive to the value of $\gamma_{2}$, approaching to infinity as $\gamma_{2} \rightarrow 1$.
For SigD, when $\gamma_{2} \rightarrow 0$, the asymptotic variance of $\lambda_{p i}$ converges to the correct asymptotic variance

$$
2 \gamma_{1}\left(s_{i}+1\right)^{2}\left(s_{i}^{2}-\gamma_{1}\right) / s_{i}^{2}
$$

of the i-th largest eigenvalue in the spiked Wishart model as derived in Paul (2007). For REG, it converges to the asymptotic variance of the i-th largest eigenvalue in the Wishart model with non-centrality spikes, derived in Onatski (2007).


Fig 2. The asymptotic variance of $\lambda_{p i}, R E G$ case. The solid and dashed lines correspond to $\gamma_{2}=0.1$ and $\gamma_{2}=0$, respectively. For both lines, $\gamma_{1}=1 / 2$.
5. Analysis of the joint density of eigenvalues. In the rest of the paper we study the statistical experiment of observing the eigenvalues of $\mathbf{F}$ when the $k$ spikes are local to some fixed points $s_{01}>\ldots>s_{0 k}$ above the phase transition threshold $\bar{s}$. The asymptotics of such an experiment can be characterized by that of the likelihood ratio corresponding to the null and alternative hypotheses

$$
H_{0}: s^{\text {true }}=s_{0} \text { and } H_{1}: s^{\text {true }}=s \equiv s_{0}+\delta / \sqrt{p}
$$

where $s_{0}=\operatorname{diag}\left\{s_{01}, \ldots, s_{0 k}\right\}$, and $\delta=\operatorname{diag}\left\{\delta_{1}, \ldots, \delta_{k}\right\}$ is the diagonal matrix of local parameters $\delta_{j} \in \mathbb{R}$. Here we instruduce notation $s^{\text {true }}$ for the true values of the spikes to contrast them with the spike parameters, $s$.

When $s^{\text {true }}=s$, the joint density of the $p$ eigenvalues of the multivariate Beta matrix $\left(n_{H} \mathbf{H}+n_{2} \mathbf{E}\right)^{-1} n_{H} \mathbf{H}$ has the following form ( see James (1964), Khatri (1967), and Muirhead (1982), pp. 312-314):

$$
\begin{equation*}
f^{\operatorname{SigD}}(x ; s)=\frac{Z_{p, n}^{\mathrm{SigD}}(x)}{\operatorname{det}\left(I_{k}+s\right)^{n_{H} / 2}}{ }_{1} F_{0}\left(N / 2 ; L s\left(I_{k}+s\right)^{-1} L^{\prime}, x\right) \tag{32}
\end{equation*}
$$

whereas for REG, we have

$$
\begin{equation*}
f^{\mathrm{REG}}(x ; s)=\frac{Z_{p, n}^{\mathrm{REG}}(x)}{\operatorname{etr}\left\{n_{H} s / 2\right\}}{ }_{1} F_{1}\left(N / 2, n_{H} / 2 ; n_{H} L s L^{\prime} / 2, x\right) \tag{33}
\end{equation*}
$$

Here the argument of the density, $x$, is a $p \times p$ real diagonal matrix; ${ }_{1} F_{0}$ and ${ }_{1} F_{1}$ are the hypergeometric functions of two matrix arguments; $N=$ $n_{H}+n_{2} ; L$ is the $p \times k$ matrix of nuisance parameters, or factor loadings, as in (2); and $Z_{p, n}^{\mathrm{Case}}(x)$ with Case $=\operatorname{SigD}$, REG depend on $n_{H}, n_{2}, p$ and $x$, but not on $s$.

Let $\tilde{\lambda}_{p j}$ be the eigenvalues of $\left(n_{H} \mathbf{H}+n_{2} \mathbf{E}\right)^{-1} n_{H} \mathbf{H}$ for some arbitrary value of $s^{\text {true }}$, not necessarily equal to $s$, and let $\tilde{\Lambda}=\operatorname{diag}\left\{\tilde{\lambda}_{p 1}, \cdots, \tilde{\lambda}_{p p}\right\}$. We would like to study the asymptotic behavior, under the null hypothesis, of the likelihood ratios

$$
f^{\text {Case }}(\tilde{\Lambda} ; s) / f^{\text {Case }}\left(\tilde{\Lambda} ; s_{0}\right)
$$

with Case $=\operatorname{SigD}, \operatorname{REG}$ as $p, n \rightarrow_{\gamma} \infty$.
The eigenvalues $\tilde{\lambda}_{p j}$ are related to the eigenvalues $\lambda_{p j}$ of the $\mathbf{F}$-ratio as follows

$$
\tilde{\lambda}_{p j}=\alpha_{n} \lambda_{p j} /\left(1+\alpha_{n} \lambda_{p j}\right), \text { where } \alpha_{n}=n_{H} / n_{2}
$$

For the purpose of the analysis of the likelihood ratios, we find it more convenient to work with $\tilde{\lambda}_{p j}$ rather than with $\lambda_{p j}$.

First, we use Lemma 1 of Passemier et al (2014) to rewrite $f^{\operatorname{SigD}}(\tilde{\Lambda} ; s)$ and $f^{\text {REG }}(\tilde{\Lambda} ; s)$ in the form of repeated contour integrals that involve hypergeometric functions of two matrix arguments of fixed dimension $k \times k$. Let $Z$ be a $k \times k$ diagonal matrix with complex variables $z_{j}$ along the diagonal, and let

$$
\begin{equation*}
\omega=\prod_{j \neq i}^{k}\left(1-z_{i} z_{j}^{-1}\right)^{1 / 2} \prod_{j=1}^{k}\left[z_{j}^{-(p-k+1) / 2} \prod_{s=1}^{p}\left(1-\tilde{\lambda}_{p s} z_{j}^{-1}\right)^{-1 / 2}\right] \tag{34}
\end{equation*}
$$

where the principal branches of all the fractional powers are taken. Let $C_{p, n}^{\text {Case }}(\tilde{\Lambda})$ be some real quantity that depend on 'Case', $n_{H}, n_{2}, p$ and $\tilde{\Lambda}$, but not on $s$; and let

$$
\begin{equation*}
k_{p, n}^{\operatorname{SigD}}(s)=\left[\operatorname{det}\left(I_{k}+s\right)\right]^{\frac{p-k-n_{H}-1}{2}}[\operatorname{det} s]^{-\frac{p-k-1}{2}} \tag{35}
\end{equation*}
$$

and

$$
\begin{equation*}
k_{p, n}^{\mathrm{REG}}(s)=\operatorname{etr}\left\{-n_{H} s / 2\right\}[\operatorname{det} s]^{-\frac{p-k-1}{2}} \tag{36}
\end{equation*}
$$

Lemma 12. Let $\tilde{\mathcal{K}}$ be a counter-clockwise oriented contour in the complex plane that encircles zero and $\tilde{\lambda}_{p j}, j=1, \ldots, p$, and intersects each of the rays $\{z: \arg z=\varphi\}, \varphi \in(-\pi, \pi]$ only once. Then, for even $p-k+1$, we have

$$
\begin{equation*}
f^{\text {Case }}(\tilde{\Lambda} ; s)=\frac{C_{p, n}^{\text {Case }}(\tilde{\Lambda}) k_{p, n}^{\text {Case }}(s)}{(2 \mathrm{i})^{k}} \int_{\tilde{\mathcal{K}}} \ldots \int_{\tilde{\mathcal{K}}} \omega \mathcal{F}_{p, n}^{\text {Case }} \prod_{i=1}^{k} \mathrm{~d} z_{i} \tag{37}
\end{equation*}
$$

where i is the imaginary unit, Case $=\operatorname{SigD}$, REG,

$$
\begin{equation*}
\mathcal{F}_{p, n}^{\mathrm{SigD}}={ }_{1} F_{0}\left(\frac{N-p+k+1}{2} ; s\left(I_{k}+s\right)^{-1}, Z\right) \tag{38}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathcal{F}_{p, n}^{\mathrm{REG}}={ }_{1} F_{1}\left(\frac{N-p+k+1}{2}, \frac{n_{H}-p+k+1}{2} ; \frac{n_{H}}{2} s, Z\right) . \tag{39}
\end{equation*}
$$

The lemma is a direct corollary of Lemma 1 of Passemier et al (2014). The requirement that $\tilde{\mathcal{K}}$ intersects each of the rays emanating from $z=0$ only once ensures that the branches of the fractional powers in $\omega$ are principal. Indeed, Onatski's (2013) Lemma 1, which Lemma 1 of Passemier et al (2014) is based on, is proven first under the assumption that $\tilde{\mathcal{K}}$ is the unit circle and the principal branches of the fractional powers in $\omega$ are used. Then the contour is deformed without changing the value of the integrals. When $\tilde{\mathcal{K}}$ is deformed so that the rays $\{z: \arg z=\varphi\}, \varphi \in(-\pi, \pi]$ are intersected by $\tilde{\mathcal{K}}$ only once, the arguments of the fractional power functions in $\omega$ never hit the negative semi-axis (note that $z_{j}^{-(p-k+1) / 2}$ is not a fractional power when $p-k+1$ is even), and therefore, the principal branches of the fractional powers should still be used after the deformation of $\tilde{\mathcal{K}}$.

In future work, it would be interesting to relax the technical requirement that $p-k+1$ is even. In a previous version of this paper, we provide such a relaxation for the case of a single spike, $k=1$. An extension to $k>1$ requires a separate non-trivial effort.
5.1. Contour deformation. Let us deform the contour of integration $\tilde{\mathcal{K}}$ into contour $\mathcal{K}$ as shown on Figure 3 . Parts $\mathcal{K}_{j}^{+}$and $\mathcal{K}_{j}^{-}, j=1, \ldots, k$, of $\mathcal{K}$ are shown non-overlaping with the real axis to enhance visibility. In fact, these parts coincide with the axis. The position $\tilde{x}_{0}$ of a kink in $\mathcal{K}$ is fixed so that

$$
\alpha \beta_{+} /\left(1+\alpha \beta_{+}\right)<\tilde{x}_{0}<\alpha x_{k} /\left(1+\alpha x_{k}\right)
$$

with $\alpha=\lim \alpha_{n}=\gamma_{2} / \gamma_{1}$, and

$$
\begin{equation*}
x_{j}=\lim l_{j}=\frac{\left(s_{0 j}+\gamma_{1}\right)\left(s_{0 j}+1\right)}{s_{0 j}-\left(s_{0 j}+1\right) \gamma_{2}}, \quad j=1, \ldots, k \tag{40}
\end{equation*}
$$

As follows from our results in the previous sections, under the null,

$$
\begin{equation*}
\tilde{\lambda}_{p k} \xrightarrow{\text { a.s. }} \alpha x_{k} /\left(1+\alpha x_{k}\right), \tag{41}
\end{equation*}
$$

and

$$
\begin{equation*}
\tilde{\lambda}_{p, k+1} \xrightarrow{\text { a.s. }} \alpha \beta_{+} /\left(1+\alpha \beta_{+}\right), \tag{42}
\end{equation*}
$$

so $\tilde{x}_{0} \in\left(\tilde{\lambda}_{p, k+1}, \tilde{\lambda}_{p k}\right)$ for sufficiently large $p$, a.s.
The radius of the circles around $\tilde{\lambda}_{p j}$ with $j=1, \ldots, k$ can be chosen arbitrarily small. Since, as can be seen from (34), the singularities of the integrand at $\tilde{\lambda}_{p j}$ are of the inverse square-root-type, the contribution of the circles to the integral disappear in the limit when the radius tends to zero. Below, we will consider this limiting version of $\mathcal{K}$, that is, the contour with the horizontal part given by the two differently oriented copies of $\left[\tilde{x}_{0}, \tilde{\lambda}_{p 1}\right]$, where the points $\tilde{\lambda}_{p 1}, \ldots, \tilde{\lambda}_{p k}$ are excluded.


Fig 3. Deformed contour $\mathcal{K}$.
Since contour $\mathcal{K}$ has common intervals with the ray $\{z: \arg z=0\}$, some of the arguments of the fractional power functions involved in $\omega$ are real and negative. Therefore, care should be taken to identify the branches used.

Suppose that

$$
z_{j_{1}}, \ldots, z_{j_{r}} \in \mathcal{K} \cap\left[\tilde{x}_{0}, \tilde{\lambda}_{p 1}\right]
$$

where $r \leq k$ and $z_{j_{1}}<\ldots<z_{j_{r}}$, and let all $z_{j}$ with $j \notin\left\{j_{1}, \ldots, j_{r}\right\}$ belong to $\mathcal{K} \backslash\left[\tilde{x}_{0}, \tilde{\lambda}_{1}\right]$. To simplify notation, we may assume that $j_{t}=t$. Since $\omega$ is symmetric in $z_{1}, \ldots, z_{k}$, this assumption is without loss of generality. Then
the parts of $\omega$ that need the branch identification are

$$
\prod_{j>i}^{r}\left(1-z_{j} z_{i}^{-1}\right)^{1 / 2} \text { and }\left(1-\tilde{\lambda}_{p s} z_{i}^{-1}\right)^{-1 / 2} \text { for } \tilde{\lambda}_{p s}>z_{i}
$$

In the SM , we prove the following lemma.
Lemma 13. Suppose that $z_{1}, \ldots, z_{r} \in \mathcal{K} \cap\left[\tilde{x}_{0}, \tilde{\lambda}_{p 1}\right]$ are such that $z_{1}<\ldots<$ $z_{r}$, and let $\operatorname{sgn}_{z_{i}}=+1$ if $z_{i}$ belongs to the "upper" portion of $\mathcal{K} \cap\left[\tilde{x}_{0}, \tilde{\lambda}_{p 1}\right]$, that is, the portion oriented from $\tilde{\lambda}_{p 1}$ to $\tilde{x}_{0}$, and $\operatorname{sgn}_{z_{i}}=-1$ if $z_{i}$ belongs to ${ }_{\tilde{\lambda}}$ the "lower" portion of $\mathcal{K} \cap\left[\tilde{x}_{0}, \tilde{\lambda}_{p 1}\right]$, that is, the portion oriented from $\tilde{x}_{0}$ to $\tilde{\lambda}_{p 1}$. Then for $j>i$, we have

$$
\left(1-z_{j} z_{i}^{-1}\right)^{1 / 2}=\mathrm{i} \times \operatorname{sgn}_{z_{i}}\left|1-z_{j} z_{i}^{-1}\right|^{1 / 2}
$$

while for $\tilde{\lambda}_{p s}>z_{i}$, we have

$$
\left(1-\tilde{\lambda}_{p s} z_{i}^{-1}\right)^{-1 / 2}=-\mathrm{i} \times \operatorname{sgn}_{z_{i}}\left|1-\tilde{\lambda}_{p s} z_{i}^{-1}\right|^{-1 / 2}
$$

5.2. Decomposition of the contour integral. Let us split $\mathcal{K}$ into $2 \times(k+1)$ parts

$$
\mathcal{K}=\bigcup_{i=1}^{k+1}\left\{\mathcal{K}_{i}^{+} \cup \mathcal{K}_{i}^{-}\right\}
$$

as shown on Figure 3, and let $\mathcal{K}_{i}=\mathcal{K}_{i}^{+} \cup \mathcal{K}_{i}^{-}$. For any $\sigma=\left(\sigma_{1}, \ldots, \sigma_{k}\right)$ with $\sigma_{i} \in\{1, \ldots, k+1\}$, let

$$
\mathcal{I}_{\sigma}=\frac{1}{(2 \mathrm{i})^{k}} \int_{\mathcal{K}_{\sigma_{k}}} \ldots \int_{\mathcal{K}_{\sigma_{1}}} \omega \mathcal{F}_{p, n} \prod_{i=1}^{k} \mathrm{~d} z_{i} .
$$

Since $\omega \mathcal{F}_{p, n}$ is symmetric in the variables $z_{i}$, we may permute them so that $z_{1}<\ldots<z_{r}$ are in $\mathcal{K} \cap\left[\tilde{x}_{0}, \tilde{\lambda}_{p 1}\right]$ and $z_{r+1}, \ldots, z_{k}$ lie in $\mathcal{K}_{k+1}$. Let $\mathcal{S}^{r}$ denote the simplex defined by $z_{1}<\ldots<z_{r}$. Consider a sequence $\sigma$ with $\sigma_{r} \leq \ldots \leq \sigma_{1} \leq k$ and $\sigma_{r+1}=\ldots=\sigma_{k}=k+1$. Consider the iterated integral

$$
\begin{equation*}
\mathcal{I}_{\sigma}^{\prime}=\frac{1}{(2 \mathrm{i})^{k}} \int_{\mathcal{K}_{\sigma} \cap \mathcal{S}^{r}} \omega \mathcal{F}_{p, n} \prod_{i=1}^{k} \mathrm{~d} z_{i} \tag{43}
\end{equation*}
$$

In the SM , we show that $\mathcal{I}_{\sigma}^{\prime}$ vanishes if there are any repeats in $\sigma$ on the real axis.

If $\mathcal{I}_{\sigma}^{\prime}$ vanishes, $\mathcal{I}_{\sigma}$ vanishes too. Therefore, the components of the integral

$$
\mathcal{I} \equiv \frac{1}{(2 \mathrm{i})^{k}} \int_{\mathcal{K}} \ldots \int_{\mathcal{K}} \omega \mathcal{F}_{p, n} \prod_{j=1}^{k} \mathrm{~d} z_{j}
$$

represented by $\mathcal{I}_{\sigma}$ with repeated $\sigma_{j} \leq k$ equal zero. This implies the following lemma. Let $\tau$ be any subset of $\{1,2, \ldots ., k\}$, and let $\sigma_{\tau}=\left(\sigma_{1 \tau}, \ldots, \sigma_{k \tau}\right)$, where

$$
\sigma_{j \tau}=\left\{\begin{array}{ll}
k+1 & \text { if } j \in \tau \\
j & \text { if } j \notin \tau
\end{array} .\right.
$$

Lemma 14. Let $T$ be the set of all the subsets of $\{1,2, \ldots ., k\}$. Then,

$$
\begin{equation*}
\mathcal{I}=\sum_{\tau \in T} \frac{k!}{|\tau|!} \mathcal{I}_{\sigma_{\tau}} \tag{44}
\end{equation*}
$$

Remark 15. The multiplier $k!/|\tau|$ ! in the latter expression counts the number of integrals $\mathcal{I}_{\sigma}$, which are different from $\mathcal{I}_{\sigma_{\tau}}$ only by permutation of the variables of integration, $z_{1}, \ldots, z_{k}$.

Below, we will show that, asymptotically as $p, n \rightarrow_{\gamma} \infty$, all integrals $\mathcal{I}_{\sigma_{\tau}}$ are dominated by

$$
\mathcal{I}_{\sigma_{\varnothing}}=\frac{1}{(2 \mathrm{i})^{k}} \int_{\mathcal{K}_{k}} \ldots \int_{\mathcal{K}_{1}} \omega \mathcal{F}_{p, n} \prod_{j=1}^{k} \mathrm{~d} z_{j}
$$

so that $\mathcal{I}$ is asymptotically equivalent to $k!\mathcal{I}_{\sigma_{\varnothing}}$. Using Lemma 13 , it is straightforward to verify that

$$
\begin{equation*}
\mathcal{I}_{\sigma_{\varnothing}}=\int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\tilde{\lambda}_{p k}}^{\tilde{\lambda}_{p, k-1}} \cdots \int_{\tilde{\lambda}_{p 2}}^{\tilde{\lambda}_{p 1}}|\omega| \mathcal{F}_{p, n} \prod_{j=1}^{k} \mathrm{~d} z_{j} \tag{45}
\end{equation*}
$$

Note that the constant $(2 \mathrm{i})^{k}$ in the denominator has canceled out.
To study the asymptotics of $\mathcal{I}_{\sigma_{\varnothing}}$, we will use Laplace approximation for the integrals involved in the above expression. However, first, we need to replace $\mathcal{F}_{p, n}$, that involves the hypergeometric function ${ }_{1} F_{0}$ for SigD and ${ }_{1} F_{1}$ for REG, by tractable approximations. This requires a separate Laplace approximation step.
5.3. Laplace approximations for $\mathcal{F}_{p, n}$. As follows from equation (38),

$$
\begin{aligned}
\mathcal{F}_{p, n}^{\mathrm{SigD}} & ={ }_{1} F_{0}\left(m ; s\left(I_{k}+s\right)^{-1}, Z\right) \\
& =[\operatorname{det}(I-Z)]^{-m}{ }_{1} F_{0}\left(m ;-\left(I_{k}+s\right)^{-1}, Z\left(I_{k}-Z\right)^{-1}\right),
\end{aligned}
$$

where $m=(N-p+k+1) / 2$. Chang (1970) studies the asymptotic behavior of ${ }_{1} F_{0}(m,-A, B)$ for fixed diagonal matrices $A$ and $B$ as $m \rightarrow \infty$. The following lemma uses a minor modification of Chang's Theorem 1 to derive a Laplace approximation for $\mathcal{F}_{p, n}^{\text {SigD }}$ that is uniform over a set of diagonal matrices $s$ and $Z$ (see the SM for a proof).

Lemma 16. For $Z=\operatorname{diag}\left\{z_{1}, \ldots, z_{k}\right\}$ such that $1>z_{1}>\ldots>z_{k}>0$, and for $s=\operatorname{diag}\left\{s_{1}, \ldots, s_{k}\right\}$ such that $s_{1}>\ldots>s_{k}>0$, as $m \rightarrow \infty$, we have

$$
\begin{equation*}
\mathcal{F}_{p, n}^{\mathrm{SigD}}=\Gamma_{k}(k / 2) \pi^{-k(k+1) / 4}\left(D^{\mathrm{SigD}}\right)^{-m} \prod_{i<j}^{k}\left(m c_{i j}\right)^{-1 / 2}(1+o(1)), \tag{46}
\end{equation*}
$$

where $\Gamma_{k}(x)=\pi^{k(k-1) / 4} \prod_{i=1}^{k} \Gamma(x-(i-1) / 2)$ is the multivariate Gamma function,

$$
D^{\mathrm{SigD}}=\operatorname{det}\left(I_{k}-s\left(I_{k}+s\right)^{-1} Z\right), c_{i j}=\frac{\left(s_{i}-s_{j}\right)\left(z_{i}-z_{j}\right)}{\left(1+s_{i}-s_{i} z_{i}\right)\left(1+s_{j}-s_{j} z_{j}\right)},
$$

and $o(1) \rightarrow 0$ uniformly on any compact subsets of the simplexes $1>z_{1}>$ $\ldots>z_{k}>0$ and $s_{1}>\ldots>s_{k}>0$.

For REG case, we have from (39)

$$
\mathcal{F}_{p, n}^{\mathrm{REG}}={ }_{1} F_{1}(T a+(k+1) / 2, T b+(k+1) / 2, T s / 2, Z),
$$

where $T=n_{H}, a=(N-p) /\left(2 n_{H}\right)$, and $b=\left(n_{H}-p\right) /\left(2 n_{H}\right)$. Note that for sufficiently large $p$, as $p, n \rightarrow_{\gamma} \infty$, we must have $a \in(1 / 2, \infty)$ and $b \in(0,1 / 2)$.

The asymptotics of ${ }_{1} F_{1}(\alpha, \beta ; A, B)$ where $\alpha$ and $A$ diverge to $\infty$ at the same rate was studied in Glynn (1980). We however need the asymptotics of this function when not only $\alpha \equiv T a+(k+1) / 2$ and $A \equiv T s / 2$, but also $\beta \equiv T b+(k+1) / 2$ diverge to infinity. Following Glynn's (1980) strategy of proof, we derive the following result. Its proof is reported in the SM.

Lemma 17. Suppose that $a$ and $b$ belong to compact subsets of $(1 / 2, \infty)$ and $(0,1 / 2)$, respectively, while the diagonal entries of $s=\operatorname{diag}\left(s_{1}, \ldots, s_{k}\right)$
and $Z=\operatorname{diag}\left(z_{1}, \ldots, z_{k}\right)$ belong to compact subsets of the simplexes $s_{1}>$ $\ldots>s_{k}>0$ and $z_{1}>\ldots>z_{k}>0$. Then, as $T \rightarrow \infty$, we have

$$
\begin{align*}
\mathcal{F}_{p, n}^{\mathrm{REG}} & =(T / 2)^{-k(k-1) / 4} \frac{\Gamma_{k}(k / 2)}{\pi^{k(k+1) / 4}} \frac{b^{k T b+k(k+1) / 4}}{a^{k T a+k(k+1) / 4}}  \tag{47}\\
& \times D^{\mathrm{REG}} \prod_{i<j}\left(z_{i}-z_{j}\right)^{-1 / 2}\left(s_{i}-s_{j}\right)^{-1 / 2} \times(1+o(1)),
\end{align*}
$$

where

$$
\begin{gathered}
D^{\mathrm{REG}}=\prod_{j=1}^{k} e^{T z_{j+}} \frac{\left(z_{j+}+a\right)^{a T}}{\left(z_{j+}+b\right)^{b T}}\left(\frac{z_{j+}\left(z_{j+}+a\right)}{z_{j+}^{2}+a \zeta_{j}}\right)^{\frac{1}{2}}, \\
z_{j+}=\frac{1}{2}\left\{\zeta_{j}-b+\sqrt{\left(b-\zeta_{j}\right)^{2}+4 a \zeta_{j}}\right\} \text { with } \zeta_{j}=z_{j} s_{j} / 2,
\end{gathered}
$$

and $o(1) \rightarrow 0$ uniformly over $a, b, s$ and $Z$ that satisfy the above requirements.
5.4. Laplace approximations for $\mathcal{I}_{\sigma_{\varnothing}}$. Note that the asymptotic approximations (46) and (47) do not hold for $z_{i}, i=1, \ldots, k$, that may approach one another. Therefore, we shall, first, analyze a multiple integral with trimmed integration domains

$$
\begin{equation*}
\overline{\mathcal{I}}_{\sigma_{\varnothing}}=\int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\tilde{\lambda}_{p k}+\varepsilon}^{\tilde{\lambda}_{p, k-1}} \cdots \int_{\tilde{\lambda}_{p 2}+\varepsilon}^{\tilde{\lambda}_{p 1}}|\omega| \mathcal{F}_{p, n} \prod_{j=1}^{k} \mathrm{~d} z_{j}, \tag{48}
\end{equation*}
$$

where $\varepsilon$ is a fixed small positive number. Then, we will show that $\overline{\mathcal{I}}_{\sigma_{\varnothing}}$ is asymptotically equivalent to $\mathcal{I}_{\sigma_{\varnothing}}$.

Although the strategy of such an analysis is the same for SigD and REG, the details are different. We start from the SigD case and then turn to the REG case.
$\underline{\text { SigD case. First, we use Lemma } 16 \text { to obtain }}$

$$
\begin{equation*}
\overline{\mathcal{I}}_{\sigma_{\varnothing}}^{\mathrm{SigD}}=b^{\mathrm{SigD}}(s) \int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\tilde{\lambda}_{p k}+\varepsilon}^{\tilde{\lambda}_{p, k-1}} \cdots \int_{\tilde{\lambda}_{p 2}+\varepsilon}^{\tilde{\lambda}_{p 1}} G^{\mathrm{SigD}} \times(1+o(1)) \prod_{j=1}^{k} \mathrm{~d} z_{j}, \tag{49}
\end{equation*}
$$

where

$$
\begin{equation*}
b^{\mathrm{SigD}}(s)=\frac{\Gamma_{k}(k / 2)}{\pi^{k(k+1) / 4} m^{k(k-1) / 4}} \prod_{i<j}^{k}\left(\frac{\left(1+s_{i}\right)\left(1+s_{j}\right)}{s_{i}-s_{j}}\right)^{1 / 2} \tag{50}
\end{equation*}
$$

and $o(1)$ converges to zero as $p, n \rightarrow_{\gamma} \infty$, uniformly over $s=s_{0}+\delta / \sqrt{p}$ such that ( $\delta_{1}, \ldots, \delta_{k}$ ) belongs to a compact subset of $\mathbb{R}^{k}$, and over $Z$ such that $\left(z_{1}, \ldots, z_{r}\right)$ belongs to the trimmed domain of integration.

Consider the inner-most integral in (49),

$$
\overline{\mathcal{I}}_{\sigma_{\varnothing}, i n}^{\mathrm{SigD}}=\int_{\tilde{\lambda}_{p 2}+\varepsilon}^{\tilde{\lambda}_{p 1}} G^{\mathrm{SigD}} \times(1+o(1)) \mathrm{d} z_{1} .
$$

Using the definition of $\omega$, we rewrite this integral in the following form

$$
\begin{equation*}
\overline{\mathcal{I}}_{\sigma_{\varnothing}, i n}^{\mathrm{SigD}}=G_{-1}^{\mathrm{SigD}} \int_{\tilde{\lambda}_{p 2}+\varepsilon}^{\tilde{\lambda}_{p 1}} e^{-p f_{1}\left(z_{1}\right)} g_{1}\left(z_{1}\right)(1+o(1)) \mathrm{d} z_{1}, \tag{51}
\end{equation*}
$$

where $f_{1}, g_{1}$, and $G_{-1}^{\mathrm{SigD}}$ are defined in Table 1 .
Table 1
Definition of $f_{1}, g_{1}$, and $G_{\mathrm{SigD},-1}$ used in equation (51), and of $H_{j}$ used in Lemma 18.

| Quantity | Definition |
| :--- | :--- |
| $f_{1}\left(z_{1}\right)$ | $=\frac{N-p+2}{2 p} \ln \left(1-\frac{s_{1}, z_{1}}{1+s_{1}}\right)+\frac{1}{2 p} \sum_{j=k+1}^{p} \ln \left(z_{1}-\lambda_{p j}\right)$ |
| $g_{1}\left(z_{1}\right)$ | $=\left(\tilde{\lambda}_{p 1}-z_{1}\right)^{-1 / 2} \prod_{j=2}^{k}\left(\left(z_{1}-z_{j}\right) /\left(z_{1}-\tilde{\lambda}_{p j}\right)\right)^{1 / 2}$ |
| $G_{-1}^{\text {SigD }}$ | $=\left\|\omega_{-1}\right\|\left(D_{-1}^{\text {SigD }}\right)^{-\frac{N-p+2}{2}} \prod_{i, j=2, \ldots, k ; ; i<j}\left(z_{i}-z_{j}\right)^{-1 / 2}$ |
| $D_{-1}^{\text {SigD }}$ | $=\operatorname{det}\left(I_{k-1}-s_{-1}\left(I_{k-1}+s_{-1}\right)^{-1} Z_{-1}\right)$ |
| $\omega_{-1}$ | $=\prod_{i, j=2, \ldots, k ; j ; \neq i}\left(1-z_{i} z_{j}^{-1}\right)^{1 / 2} \prod_{j=2}^{k}\left[z_{j}^{-(p-k+2) / 2} \prod_{s=1}^{p}\left(1-\tilde{\lambda}_{p s} z_{j}^{-1}\right)^{-1 / 2}\right]$ |
| $Z_{-1}$ | $=\operatorname{diag}\left\{z_{2}, \ldots, z_{k}\right\}$ |
| $s_{-1}$ | $=\operatorname{diag}\left\{s_{2}, \ldots, s_{k}\right\}$ |
| $H_{j}$ | $=H_{j, \text { num }} / H_{j, \text { den }}$ |
| $H_{j, \text { num }}$ | $=\left(1-\gamma_{2}\right) s_{0 j}\left(\gamma_{1}+\gamma_{2}+\gamma_{2} s_{0 j}\right)\left(1+s_{0 j}-\beta_{+}^{1 / 2}\right)\left(1+s_{0 j}-\beta_{-}^{1 / 2}\right)$ |
| $H_{j, \text { den }}$ | $=2 \gamma_{1} \gamma_{2}\left(s_{0 j}-\gamma_{2} s_{0 j}-\gamma_{2}\right)\left(1+s_{0 j}\right)\left(\gamma_{1}+s_{0 j}\right)$ |

Using the Laplace method to approximate the integral in (51) (see the SM for details), and then repeating the procedure for the second, third, etc. to the inner-most integral in (49) and combining the results, we obtain the following lemma.

Lemma 18. Under the null hypothesis $H_{0}: s^{\text {true }}=s_{0}$, as $p, n \rightarrow_{\gamma} \infty$, $\overline{\mathcal{I}}_{\sigma_{\varnothing}}^{\mathrm{SigD}}=b^{\mathrm{SigD}}(s) \prod_{j=1}^{k}\left[\Omega_{j}^{\mathrm{SigD}}\left(p H_{j} / \pi\right)^{-1 / 2} \prod_{s=k+1}^{p}\left(\tilde{\lambda}_{p j}-\tilde{\lambda}_{p s}\right)^{-1 / 2}\right](1+o(1))$,
where $\Omega_{j}^{\mathrm{SigD}}=\left(1-\frac{s_{j}}{1+s_{j}} \tilde{\lambda}_{p j}\right)^{-\frac{N-p+2}{2}}, H_{j}$ are as defined in Table 1, and $o(1) \rightarrow 0$ uniformly over $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, a.s.
$\underline{R E G \text { case. First, we use Lemma 17, to obtain }}$

$$
\begin{equation*}
\overline{\mathcal{I}}_{\sigma_{\varnothing}}^{\mathrm{REG}}=b^{R E G}(s) \int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\tilde{\lambda}_{p k}+\varepsilon}^{\tilde{\lambda}_{p, k-1}} \cdots \int_{\tilde{\lambda}_{p 2}+\varepsilon}^{\tilde{\lambda}_{p 1}} G^{\mathrm{REG}} \times(1+o(1)) \prod_{j=1}^{k} \mathrm{~d} z_{j} \tag{52}
\end{equation*}
$$

where

$$
G^{\mathrm{REG}}=|\omega| D^{\mathrm{REG}} \prod_{i<j}\left(z_{i}-z_{j}\right)^{-1 / 2}
$$

with

$$
\begin{gather*}
D^{\mathrm{REG}}=\prod_{j=1}^{k} e^{n_{H} z_{j+}} \frac{\left(z_{j+}+a\right)^{a n_{H}}}{\left(z_{j+}+b\right)^{b n_{H}}}\left(\frac{z_{j+}\left(z_{j+}+a\right)}{z_{j+}^{2}+a \zeta_{j}}\right)^{\frac{1}{2}} \\
a=(N-p) /\left(2 n_{H}\right), b=\left(n_{H}-p\right) /\left(2 n_{H}\right), \text { and } \\
b^{\mathrm{REG}}(s)=\left(\frac{n_{H}}{2}\right)^{-k(k-1) / 4} \frac{\Gamma_{k}(k / 2)}{\pi^{k(k+1) / 4}} \frac{b^{k n_{H} b+k(k+1) / 4}}{a^{k n_{H} a+k(k+1) / 4}} \prod_{i<j}\left(s_{i}-s_{j}\right)^{-1 / 2} \tag{53}
\end{gather*}
$$

The same uniformity properties of $o(1)$ as in the case of (49) apply.
Consider the inner-most integrals in (52),

$$
\overline{\mathcal{I}}_{\sigma_{\varnothing}, i n}^{\mathrm{REG}}=\int_{\tilde{\lambda}_{p 2}+\varepsilon}^{\tilde{\lambda}_{p 1}} G^{\mathrm{REG}} \times(1+o(1)) \mathrm{d} z_{1}
$$

Using the definition of $\omega$, we rewrite this integral in the following form

$$
\begin{equation*}
\overline{\mathcal{I}}_{\sigma_{\varnothing}, i n}^{\mathrm{REG}}=G_{-1}^{\mathrm{REG}} \int_{\tilde{\lambda}_{p 2}+\varepsilon}^{\tilde{\lambda}_{p 1}} e^{-n_{H} f_{2}\left(z_{1}\right)} g_{2}\left(z_{1}\right)(1+o(1)) \mathrm{d} z_{1} \tag{54}
\end{equation*}
$$

where $f_{2}, g_{2}$, and $G_{-1}^{\mathrm{REG}}$ are defined in Table 2.
Similarly to the integral in (51), the one in (54) can be analyzed using the Laplace approximation steps. Repeating the procedure for the second, third, etc. to the inner-most integral in (52) and combining the results, we obtain the following lemma. See the SM for a detailed proof.

Lemma 19. Under the null hypothesis $H_{0}: s^{\text {true }}=s_{0}$, as $p, n \rightarrow \gamma \infty$, $\overline{\mathcal{I}}_{\sigma_{\varnothing}}^{\mathrm{REG}}=b^{\mathrm{REG}}(s) \prod_{j=1}^{k}\left[\Omega_{j}^{\mathrm{REG}}\left(p H_{j} / \pi\right)^{-1 / 2} \prod_{s=k+1}^{p}\left(\tilde{\lambda}_{p j}-\tilde{\lambda}_{p s}\right)^{-1 / 2}\right](1+o(1))$,
where

$$
\Omega_{j}^{\mathrm{REG}}=e^{n_{H} \bar{z}_{j+}} \frac{\left(\bar{z}_{j+}+a\right)^{a n_{H}}}{\left(\bar{z}_{j+}+b\right)^{b n_{H}}} \frac{\gamma_{1}+\gamma_{2}+\gamma_{2} s_{0 j}}{\left(\left(\gamma_{1}+\gamma_{2}+\gamma_{2} s_{0 j}\right)^{2}-\rho^{2} \gamma_{1}\right)^{1 / 2}}
$$

and $\bar{z}_{j+}$ is the value of $z_{j+}$ that corresponds to $\zeta_{j}=\tilde{\lambda}_{p j} s_{j} / 2$. The $H_{j}$ are as defined in Table 1, and $o(1) \rightarrow 0$ uniformly over $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, a.s.

Now let us show that $\overline{\mathcal{I}}_{\sigma_{\varnothing}}$ is asymptotically equivalent to $\mathcal{I}_{\sigma_{\varnothing}}$. By definition,

$$
\mathcal{I}_{\sigma_{\varnothing}}-\overline{\mathcal{I}}_{\sigma_{\varnothing}}=\sum_{\mathcal{D}} \int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\mathcal{D}_{k-1}} \ldots \int_{\mathcal{D}_{1}}|\omega| \mathcal{F}_{p, n} \prod_{j=1}^{k} \mathrm{~d} z_{j},
$$

where the sum runs over all $\mathcal{D}_{j}$ that are represented by either $\left[\tilde{\lambda}_{j+1}+\varepsilon, \tilde{\lambda}_{j}\right]$ or $\left[\tilde{\lambda}_{j+1}, \tilde{\lambda}_{j+1}+\varepsilon\right]$, and at least one $\mathcal{D}_{j}, j=1, \ldots, k-1$, is represented by $\left[\tilde{\lambda}_{j+1}, \tilde{\lambda}_{j+1}+\varepsilon\right]$. All terms in this sum cam be analyzed similarly. Let us explain the main idea of the analysis using the term

$$
J \equiv \int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\tilde{\lambda}_{p k}+\varepsilon}^{\tilde{\lambda}_{p, k-1}} \cdots \int_{\tilde{\lambda}_{p 3}+\varepsilon}^{\tilde{\lambda}_{p 2}} \int_{\tilde{\lambda}_{p 2}}^{\tilde{\lambda}_{p 2}+\varepsilon}|\omega| \mathcal{F}_{p, n} \prod_{j=1}^{k} \mathrm{~d} z_{j} .
$$

Since the lower integration limit of the inner-most integral coincides with the upper integration limit of the second inner-most integral, we cannot

Table 2
Definition of $f_{2}, g_{2}$, and $G_{\text {REG },-1}$ used in equation (54). Quantities $\omega_{-1}, Z_{-1}$, and $s_{-1}$ are as defined in Table 1.

| Quantity | Definition |
| :--- | :--- |
| $f_{2}\left(z_{1}\right)$ | $=-z_{1+}-a \ln \left(z_{1+}+a\right)+b \ln \left(z_{1+}+b\right)+\frac{1}{2 n_{H}} \sum_{j=k+1}^{p} \ln \left(z_{1}-\lambda_{p j}\right)$ |
| $g_{2}\left(z_{1}\right)$ | $=\left(\frac{z_{1}+\left(z_{1+}+a\right)}{z_{1}^{2}+a \zeta_{1}}\right)^{\frac{1}{2}}\left(\tilde{\lambda}_{p 1}-z_{1}\right)^{-1 / 2} \prod_{j=2}^{k}\left(\frac{z_{1}-z_{j}}{z_{1} \lambda_{p j}}\right)^{1 / 2}$ |
| $G_{-1}^{\mathrm{REG}}$ | $=\left\|\omega_{-1}\right\| D_{-1}^{\mathrm{REG}} \prod_{i, j=2, \ldots, k ; i<j}\left(z_{i}-z_{j}\right)^{-1 / 2}$ |
| $D_{-1}^{\mathrm{REG}}$ | $=\prod_{j=2}^{k} e^{n_{H} z_{j}+\frac{\left(z_{j+}+a\right)^{a n_{H}}}{\left(z_{j+}+b\right)^{b_{H}}}\left(\frac{z_{j+}\left(z_{j++}\right)}{z_{j+}^{2}+a \zeta_{j}}\right)^{\frac{1}{2}}}$ |

use Lemmas 16 and 17 to approximate $\mathcal{F}_{p, n}$ uniformly over the integration domain of $J$. However, we can obtain an upper bound on $|J|$ that can be analyzed using these lemmas.

The key is to observe that $\mathcal{F}_{p, n}$ viewed as a function of $Z \equiv \operatorname{diag}\left\{z_{1}, \ldots, z_{k}\right\}$ is positive and monotonically increasing in each of $0<z_{j}<1, j=1, \ldots, k$, for all $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$ and all sufficiently large $p$. This follows from the representation of ${ }_{1} F_{0}$ and ${ }_{1} F_{1}$ in the series of zonal polynomials and from the monotonicity of the zonal polynomials of $Z$ in each of $0<z_{j}<1, j=1, \ldots, k$. Such a monotonicity follows from the fact that zonal polynomials are linear combinations of monomial symmetric functions of $z_{j}$ with positive coefficients (see Chattopahyay and Pillai (1970), Lemma 2).

Let us make the dependence of $\mathcal{F}_{p, n}$ on $Z$ explicit by writing $\mathcal{F}_{p, n}(Z)$. The positivity and monotonicity of $\mathcal{F}_{p, n}(Z)$ yield the following bound

$$
\begin{equation*}
|J| \leq \int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\tilde{\lambda}_{p k}+\varepsilon}^{\tilde{\lambda}_{p, k-1}} \cdots \int_{\tilde{\lambda}_{p 3}+\varepsilon}^{\tilde{\lambda}_{p 2}} \int_{\tilde{\lambda}_{p 2}}^{\tilde{\lambda}_{p 2}+\varepsilon}|\omega| \mathcal{F}_{p, n}\left(\tilde{Z}_{1}\right) \prod_{j=1}^{k} \mathrm{~d} z_{j} \tag{55}
\end{equation*}
$$

where $\tilde{Z}_{1}=\operatorname{diag}\left\{\tilde{\lambda}_{p 2}+\varepsilon, z_{2}, \ldots, z_{k}\right\}$.
In constrast to $\mathcal{F}_{p, n}(Z)$, function $\mathcal{F}_{p, n}\left(\tilde{Z}_{1}\right)$ can be approximated using Lemmas 16 and 17. Exploiting such an approximation to show that the right hand side of (55) is asymptotically dominated by $\overline{\mathcal{I}}_{\sigma_{\varnothing}}$ yields the following lemma. Its proof is given in the SM.

Lemma 20. Under the null hypothesis $H_{0}: s^{\text {true }}=s_{0}$, as $p, n \rightarrow_{\gamma} \infty$,

$$
\mathcal{I}_{\sigma_{\varnothing}}=\overline{\mathcal{I}}_{\sigma_{\varnothing}}(1+o(1)),
$$

where $o(1) \rightarrow 0$ uniformly over $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, a.s.
5.5. Asymptotic negligibility of $\mathcal{I}_{\sigma_{\tau}}$ with $\tau \neq \varnothing$, and a summary. To finish our asymptotic analysis of the joint densities of the eigenvalues, we need to show that integrals $\mathcal{I}_{\sigma_{\tau}}$ with $\tau \neq \varnothing$ are asymptotically dominated by $\mathcal{I}_{\sigma_{\varnothing}}$. This can be established similarly to Lemma 20.

Consider, for example, $\tau=\{1,2, \ldots, k\}$. By definition,

$$
\mathcal{I}_{\sigma_{\tau}}=\frac{1}{(2 \mathrm{i})^{k}} \int_{\mathcal{K}_{k+1}} \ldots \int_{\mathcal{K}_{k+1}} \omega \mathcal{F}_{p, n}(Z) \prod_{j=1}^{k} \mathrm{~d} z_{j}
$$

Since $\mathcal{K}_{k+1}$ is not a subset of $\mathbb{R}$, we cannot use Lemmas 16 and 17 to approximate $\mathcal{F}_{p, n}(Z)$. However, it is easy to obtain an upper bound on $\left|\mathcal{F}_{p, n}(Z)\right|$ that can be approximated using those lemmas.

Indeed, the representation of ${ }_{1} F_{0}$ and ${ }_{1} F_{1}$ in the series of zonal polynomials and the fact that these polynomials are linear combinations of monomial symmetric functions with positive coefficients (see Chattopahyay and Pillai (1970), Lemma 2) yield the following inequality

$$
\left|\mathcal{F}_{p, n}(Z)\right| \leq \mathcal{F}_{p, n}(|Z|)
$$

where $|Z|=\operatorname{diag}\left\{\left|z_{1}\right|, \ldots,\left|z_{k}\right|\right\}$. Therefore, we have

$$
\left|\mathcal{I}_{\sigma_{\tau}}\right| \leq \frac{1}{2^{k}} \int_{\mathcal{K}_{k+1}} \ldots \int_{\mathcal{K}_{k+1}}|\omega| \mathcal{F}_{p, n}\left(\tilde{x}_{0} I_{k}\right) \prod_{j=1}^{k}\left|\mathrm{~d} z_{j}\right| .
$$

Further, by the monotonicity of zonal polynomials, $\mathcal{F}_{p, n}\left(\tilde{x}_{0} I_{k}\right) \leq \mathcal{F}_{p, n}\left(Z_{\eta}\right)$, where $Z_{\eta}=\operatorname{diag}\left\{\tilde{x}_{0}+k \eta, \ldots, \tilde{x}_{0}+2 \eta, \tilde{x}_{0}+\eta\right\}$ and $\eta$ is a fixed small positive number. Therefore,

$$
\left|\mathcal{I}_{\sigma_{\tau}}\right| \leq \frac{1}{2^{k}} \int_{\mathcal{K}_{k+1}} \ldots \int_{\mathcal{K}_{k+1}}|\omega| \mathcal{F}_{p, n}\left(Z_{\eta}\right) \prod_{j=1}^{k}\left|\mathrm{~d} z_{j}\right|
$$

Function $\mathcal{F}_{p, n}\left(Z_{\eta}\right)$ can now be approximated using Lemmas 16 and 17 , which yields the following lemma (see the SM for a proof).

Lemma 21. Under the null hypothesis $H_{0}: s^{\text {true }}=s_{0}$, as $p, n \rightarrow \gamma \infty$, for any $\tau \neq \varnothing$,

$$
\mathcal{I}_{\sigma_{\tau}}=o(1) \mathcal{I}_{\sigma_{\varnothing}}
$$

where $o(1) \rightarrow 0$ uniformly over $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, a.s.

In conclusion of this section, we formulate a theorem that describes the asymptotic behavior of the joint density of the eigenvalues of the multivariate Beta matrix $\left(n_{H} \mathbf{H}+n_{2} \mathbf{E}\right)^{-1} n_{H} \mathbf{H}$ by combining results of Lemma 12 with those of Lemmas 18-21. For the reader's convenience, we reproduce the definitions of the quantities used in the statement of the theorem in Table 3.

Theorem 22. Under the null hypothesis $H_{0}: s^{\text {true }}=s_{0}$, as $p, n \rightarrow_{\gamma} \infty$ while $p-k+1$ remains even,

$$
\begin{aligned}
f^{\text {Case }}(\tilde{\Lambda} ; s)= & k!C_{p, n}^{\mathrm{Case}}(\tilde{\Lambda}) k_{p, n}^{\mathrm{Case}}(s) b^{\text {Case }}(s) \\
& \times \prod_{j=1}^{k}\left[\Omega_{j}^{\text {Case }}\left(p H_{j} / \pi\right)^{-1 / 2} \prod_{s=k+1}^{p}\left(\tilde{\lambda}_{p j}-\tilde{\lambda}_{p s}\right)^{-1 / 2}\right](1+o(1)),
\end{aligned}
$$

where $o(1) \rightarrow 0$ uniformly over $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, a.s.

Table 3
Definitions of $C_{p, n}^{\text {Case }}(\tilde{\Lambda}), k_{p, n}^{\text {Case }}(s), b^{\text {Case }}(s), \Omega_{j}^{\text {Case }}$, and $H_{j}$.

| Quantity | Definition |
| :--- | :--- |
| $C_{p, n}^{\text {Case }}(\tilde{\Lambda})$ | $=$ A quantity that depends on 'Case', $n_{H}, n_{2}, p$ and $\tilde{\Lambda}$, but not on $s$ |
| $k_{p, n}^{\text {SigD }}(s)$ | $=\left[\operatorname{det}\left(I_{k}+s\right)\right]^{\frac{p-k-n_{H}-1}{2}}[\operatorname{det} s]^{-\frac{p-k-1}{2}}$, |
| $k_{p, n}^{\mathrm{REG}}(s)$ | $=\operatorname{etr}\left\{-n_{H} s / 2\right\}[\operatorname{det} s]^{-\frac{p-k-1}{2}}$, |
| $b^{\mathrm{SigD}}(s)$ | $=\frac{\Gamma_{k}(k / 2)}{\pi^{k(k+1) / 4 m^{k(k-1) / 4}} \prod_{i<j}^{k}\left(\frac{\left(1+s_{i}\right)\left(1+s_{j}\right)}{s_{i}-j_{j}}\right)^{1 / 2},}$, |
| $b^{\mathrm{REG}}(s)$ | $=\left(\frac{n_{H}}{2}\right)^{-k(k-1) / 4} \frac{\Gamma_{k}(k / 2)}{\pi^{k(k+1) / 4} \frac{b^{k n_{H} b+k(k+1) / 4}}{a^{k n_{H} a+k(k+1) / 4}} \prod_{i<j}\left(s_{i}-s_{j}\right)^{-1 / 2},}$ |
| $\Omega_{j}^{\text {SigD }}$ | $=\left(1-\frac{s_{j}}{1+s_{j}} \tilde{\lambda}_{p j}\right)^{-\frac{N-p+2}{2}}$, |
| $\Omega_{j}^{\mathrm{REG}}$ | $=e^{n_{H} \bar{z}_{j}+\frac{\left(\bar{z}_{j+}+a\right)^{a n_{H}}}{\left(\bar{z}_{j+}+b\right)^{b n_{H}}} \frac{\gamma_{1}+\gamma_{2}+\gamma_{2} s_{0 j}}{\left(\left(\gamma_{1}+\gamma_{2}+\gamma_{2} s_{0 j}\right)^{2}-\rho^{2} \gamma_{1}\right)^{1 / 2}},}$ |
| $H_{j}$ | $=H_{j, \text { num }} / H_{j, \text { den }}$ |
| $H_{j, \text { num }}$ | $=\left(1-\gamma_{2}\right) s_{0 j}\left(\gamma_{1}+\gamma_{2}+\gamma_{2} s_{0 j}\right)\left(1+s_{0 j}-\beta_{+}^{1 / 2}\right)\left(1+s_{0 j}-\beta_{-}^{1 / 2}\right)$ |
| $H_{j, \text { den }}$ | $=2 \gamma_{1} \gamma_{2}\left(s_{0 j}-\gamma_{2} s_{0 j}-\gamma_{2}\right)\left(1+s_{0 j}\right)\left(\gamma_{1}+s_{0 j}\right)$ |

6. Local Asymptotic Normality. Our goal is to understand the asymptotic behavior of the likelihood ratios for the eigenvalues of the $\mathbf{F}$-ratio (or, equivalently, of the multivatiate Beta) at local alternatives to the null of $k$ supercritical spikes:

$$
H_{0}: s^{\text {true }}=s_{0}, H_{1}: s^{\text {true }}=s_{0}+\delta / \sqrt{p}
$$

Let us, first, reparametrize the alternative by considering new local parameters

$$
\theta_{j}=\delta_{j} / w\left(s_{0 j}\right) \text { for } j=1, \ldots, k
$$

where $w\left(s_{0 j}\right)$ is the component of the asymptotic variance of $\lambda_{p j}$ defined in Theorem 10. That is,

$$
w\left(s_{0 j}\right)= \begin{cases}2\left(\rho s_{0 j}\left(1+s_{0 j}\right) / d\left(s_{0 j}\right)\right)^{2} & \text { for SigD } \\ 2\left(\rho s_{0 j}\left(1+s_{0 j}\right) / d\left(s_{0 j}\right)\right)^{2}-2 \gamma_{1} s_{0 j}^{2} t\left(s_{0 j}\right) & \text { for REG }\end{cases}
$$

with

$$
\begin{aligned}
t\left(s_{0 j}\right) & =\left(1-\gamma_{2}\right)\left(s_{0 j}-\underline{\mathbf{s}}\right)\left(s_{0 j}-\bar{s}\right) / d^{2}\left(s_{0 j}\right) \text { and } \\
d\left(s_{0 j}\right) & =\left(1-\gamma_{2}\right) s_{0 j}-\gamma_{2}
\end{aligned}
$$

Let $\Lambda=\operatorname{diag}\left\{\lambda_{p 1}, \ldots, \lambda_{p p}\right\}$. Denote the likelihood ratio as

$$
L_{p}(\theta, \Lambda) \equiv f(\tilde{\Lambda} ; s) / f\left(\tilde{\Lambda} ; s_{0}\right)
$$

We write $L_{p}(\theta, \Lambda)$ instead of $L_{p}(\theta, \tilde{\Lambda})$ to emphasize the fact that the likelihood ratio remains the same whether we define it as the ratio of the joint densities of the eigenvalues of the multivariate Beta $\left(n_{H} \mathbf{H}+n_{2} \mathbf{E}\right)^{-1} n_{H} \mathbf{H}$ (the diagonal elements of $\tilde{\Lambda}$ ) or of the eigenvalues of the $\mathbf{F}$-ratio (the diagonal elements of $\Lambda$ ).

Using Theorem 22 to express $\ln L_{p}(\theta, \Lambda)$ in terms of elementary functions of $\theta$, expanding the result in the powers of $\theta p^{-1 / 2}$ up to and including terms with $\theta^{2} p^{-1}$, and invoking Theorem 4 yields the following theorem. Its proof can be found in the SM.

Theorem 23. Under the null hypothesis $H_{0}: s^{\text {true }}=s_{0}$, as $p, n \rightarrow \gamma \infty$ while $p-k+1$ remains even,

$$
\ln L_{p}(\theta, \Lambda)=\sum_{j=1}^{k}\left\{\theta_{j} \sqrt{p}\left(\lambda_{p j}-l_{j}\right)-\frac{1}{2} \theta_{j}^{2} \tau^{2}\left(s_{0 j}\right)\right\}+o_{\mathrm{P}}(1)
$$

where $l_{j}=\left(s_{0 j}+c_{1}\right)\left(s_{0 j}+1\right) /\left(\left(1-c_{2}\right) s_{0 j}-c_{2}\right), \tau^{2}\left(s_{0 j}\right)=w\left(s_{0 j}\right) t\left(s_{0 j}\right)$, and $o_{\mathrm{P}}(1) \rightarrow 0$ in probability, uniformly in $\left(\theta_{1}, \ldots, \theta_{k}\right)$ from any compact subset of $\mathbb{R}^{k}$.

Theorem 23 together with the joint asymptotic normality of $\sqrt{p}\left(\lambda_{p j}-l_{j}\right)$, $j=1, \ldots, k$, established in Theorem 10 imply, via Le Cam's First Lemma (see van der Vaart (1998), p.88), that the sequences of the probability measures $\left\{\mathbb{P}_{s_{0}, p}\right\}$ and $\left\{\mathbb{P}_{s, p}\right\}$ describing the joint distribution of the eigenvalues of $\mathbf{F}$ under the null $H_{0}: s^{\text {true }}=s_{0}$ and under the local alternative $H_{1}: s^{\text {true }}=$ $s_{0}+\theta w / \sqrt{p}$, where

$$
\theta=\operatorname{diag}\left\{\theta_{1}, \ldots, \theta_{k}\right\} \text { and } w=\operatorname{diag}\left\{w\left(s_{01}\right), \ldots, w\left(s_{0 k}\right)\right\},
$$

are mutually contiguous. Moreover, the experiments

$$
\mathcal{E}_{\theta, s_{0}, p} \equiv\left(\left(\lambda_{p 1}, \ldots, \lambda_{p p}\right) \sim \mathbb{P}_{s_{0}+\theta w / \sqrt{\bar{p}}, p}: \theta \in \mathbb{R}^{k}\right)
$$

converge to the Gaussian shift experiment

$$
\mathcal{E}_{\theta, s_{0}} \equiv\left(Y \sim N(\mu(\theta), \mathcal{T}): \theta \in \mathbb{R}^{k}\right),
$$

where

$$
\begin{aligned}
Y & =\sqrt{p}\left(\lambda_{p 1}-l_{1}, \ldots, \lambda_{p k}-l_{k}\right)^{\prime}, \\
\mu(\theta) & =\left(\theta_{1} \tau^{2}\left(s_{01}\right), \ldots, \theta_{k} \tau^{2}\left(s_{0 k}\right)\right), \text { and } \\
\mathcal{T} & =\operatorname{diag}\left\{\tau^{2}\left(s_{01}\right), \ldots, \tau^{2}\left(s_{0 k}\right)\right\} .
\end{aligned}
$$

In particular, these experiments are $L A N$.
As discussed in the introduction, the LAN property of the experiments $\mathcal{E}_{\theta, s_{0}, p}$ imply that the asymptotically efficient tests of hypotheses about super-critical spikes are based on $\lambda_{p 1}, \ldots, \lambda_{p k}$. Such tests may ignore information contained in the other eigenvalues of the $\mathbf{F}$-ratio.

Here, we will illustrate the $L A N$ property by constructing $L A N$ confidence sets for $s$. The likelihood ratio confidence set, $C S_{L R}$, is the set of all $s$ that are not rejected by the likelihood ratio test. Asymptotically, this will coincide with the set of all $s_{0}$ such that the hypothesis $H_{0}: \theta=0$ is not rejected in the limiting experiement $\mathcal{E}_{\theta, s_{0}}$. Therefore, we find the asymptotic $100(1-\alpha) \%$ $C S_{L R}$ for supercritical spikes, by collecting all $s_{0}$ that satisfy the inequality

$$
\sum_{j=1}^{k} \frac{p\left(\lambda_{p j}-l_{j}\left(s_{0 j}\right)\right)^{2}}{\hat{\tau}^{2}\left(s_{0 j}\right)} \leq \chi_{\alpha}^{2}(k),
$$

where $\chi_{\alpha}^{2}(k)$ is the critical value of the chi-squard distribution with $k$ degrees of freedom and the other quantities are as defined in Table 4.

Table 4
Definitions of $l(x), \hat{\tau}(x)$ and related quantities.

| Quantity | Definition |
| :--- | :--- |
| $l(x)$ | $=\left(x+c_{1}\right)(x+1) /\left(\left(1-c_{2}\right) x-c_{2}\right)$, |
| $\hat{\tau}^{2}(x)$ | $=\hat{w}(x) \hat{t}(x)$, |
| $\hat{w}(x)$ | $=2(r x(1+x) / \hat{d}(x))^{2}$ for SigD, |
| $\hat{w}(x)$ | $=2(r x(1+x) / \hat{d}(x))^{2}-2 c_{1} x^{2} \hat{t}(x)$ for REG, |
| $\hat{t}(x)$ | $=\left(1-c_{2}\right)(x-\hat{\hat{s}})(x-\hat{\bar{s}}) / \hat{d}^{2}(x),$, |
| $\hat{d}(x)$ | $=\left(1-c_{2}\right) x-c_{2}$, |
| $\underline{\hat{s}}$ | $=-\left(c_{1}+c_{2}\right) /\left(c_{2}+r\right)$, |
| $\hat{\bar{s}}$ | $=\left(c_{2}+r\right) /\left(1-c_{2}\right)$, |
| $r^{2}$ | $=c_{1}+c_{2}-c_{1} c_{2}$. |

Figure 4 shows the $95 \%$ asymptotic confidence sets for ( $s_{1}, s_{2}$ ) when $c_{1}=c_{2}=1 / 2, p=100$, and $\lambda_{p 1}=45, \lambda_{p 2}=25$. The outer and inner
ovals represent the confidence sets for SigD and REG cases, respectively. It is worth noting that the asymptotic confidence sets do not necessarily preserve the ranking $s_{1}>s_{2}$. Indeed, in the figure, the confidence set for the SigD case intersects the 45 -degree (dashed) line. Of course, this undesirable phenomenon will not be observed for sufficiently large $p$.


Fig 4. Confidence sets.

To assess the quality of the $L A N$ confidence sets, we conduct a small-scale Monte Carlo experiment. Specifically, we generate 10,000 replications of $\mathbf{H}$ and $\mathbf{E}$, distributed as

$$
\left(n_{1}+k\right) \mathbf{H} \sim W_{p}\left(n_{1}+k, \Sigma_{1}, \Omega_{1}\right) \quad \text { and } \quad n_{2} \mathbf{E} \sim W_{p}\left(n_{2}, I\right)
$$

and compute the eigenvalues of the corresponding $\mathbf{F}$-ratio, $\mathbf{E}^{-1} \mathbf{H}$. We consider various values of $n_{1}, n_{2}$, and $p$, and several different values of the spikes.

The Monte Carlo coverage rates of the nominal $95 \%$ confidence sets are reported in Table 5 for $k=1$ and Table 6 for $k=2$.

Table 5
Coverage probabilities, nominal $95 \%$ confidence sets. Single spike.

| SigD | REG |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | $s=3$ | $s=5$ | $s=10$ | $s=3$ | $s=5$ | $s=10$ |
| $c_{2}=0, n_{1}=p=100$ | $\hat{\hat{s}}=1.0$ | 95.8 | 95.1 | 95.5 | 95.1 | 94.9 | 95.3 |
| $n_{1}=n_{2}=100, p=50$ | $\hat{\bar{s}}=2.7$ | 79.1 | 94.2 | 94.3 | 79.6 | 94.4 | 94.3 |
| $n_{1}=n_{2}=100, p=5$ | $\hat{\bar{s}}=0.4$ | 94.9 | 94.9 | 95.1 | 94.2 | 95.0 | 94.6 |
| $n_{1}=n_{2}=100, p=2$ | $\hat{\bar{s}}=0.2$ | 95.0 | 95.3 | 95.3 | 94.9 | 95.1 | 94.5 |

Table 6
Coverage probabilities, nominal $95 \%$ confidence sets. Two spikes.

| $k=2$ |  |  | SigD |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | $s_{1}=4$ | $s_{1}=5$ | $s_{1}=10$ | $s_{1}=6$ | $s_{1}=20$ |
|  | $s_{2}=3$ | $s_{2}=3$ | $s_{2}=3$ | $s_{2}=5$ | $s_{2}=10$ |  |
| $c_{2}=0, n_{1}=p=100$ | $\hat{\bar{s}}=1.0$ | 95.9 | 95.9 | 95.6 | 95.4 | 95.6 |
| $n_{1}=n_{2}=100, p=50$ | $\hat{\bar{s}}=2.7$ | 68.1 | 75.4 | 80.8 | 92.3 | 94.9 |
| $n_{1}=n_{2}=100, p=5$ | $\hat{\bar{s}}=0.4$ | 94.8 | 95.3 | 94.3 | 94.1 | 95.2 |
| $n_{1}=n_{2}=100, p=2$ | $\hat{\bar{s}}=0.2$ | 94.7 | 95.4 | 95.1 | 94.0 | 94.9 |
|  |  |  |  | REG |  |  |
| $k=2$ |  | $s_{1}=4$ | $s_{1}=5$ | $s_{1}=10$ | $s_{1}=6$ | $s_{1}=20$ |
|  |  | $s_{2}=3$ | $s_{2}=3$ | $s_{2}=3$ | $s_{2}=5$ | $s_{2}=10$ |
| $c_{2}=0, n_{1}=p=100$ | $\hat{\bar{s}}=1.0$ | 95.5 | 95.5 | 95.3 | 95.5 | 95.1 |
| $n_{1}=n_{2}=100, p=50$ | $\hat{\bar{s}}=2.7$ | 70.4 | 76.1 | 80.5 | 92.2 | 94.2 |
| $n_{1}=n_{2}=100, p=5$ | $\hat{\bar{s}}=0.4$ | 94.8 | 95.5 | 95.1 | 93.9 | 94.5 |
| $n_{1}=n_{2}=100, p=2$ | $\hat{\bar{s}}=0.2$ | 95.0 | 95.4 | 94.4 | 94.0 | 94.4 |

Surprisingly, for $k=1$, the coverage rate of the nominal $95 \%$ confidence intervals remains good even for an extremely small dimensionality $-p=2$. The worst reported coverage corresponds to $n_{1}=n_{2}=100, p=50$ and $s=3$. This is an example of a situation where the true spike is close to the phase transition threshold, given here by $\hat{\bar{s}} \equiv\left(c_{2}+r\right) /\left(1-c_{2}\right) \approx 2.7$. A further analysis suggests that the reason for the poor coverage in this
situation is that the finite sample variance of $\lambda_{p 1}$ is substantially larger than its asymptotic counterpart.

For $k=2$, the results are similar. Again, surprisingly, the coverage remains good even when $p=2$. The worst results correspond to situations where a spike is close to the phase transition. A particularly unfavourable situation arises when both spikes are close to the threshold $-s_{1}=4, s_{2}=3$ with $\hat{s}=2.7$. A more detailed analysis shows that in such a case the asymptotic variances are smaller than the finite sample ones. In addition, the smallest of the spikes tends to lie below the corresponding a.s. limit, whereas the largest one tends to lie above the corresponding a.s. limit. These deviations become smaller when at least one of the spikes is moved away from the phase transition threshold.
7. Conclusion. In this paper, we establish the Local Asymptotic Normality of the experiments of observing the eigenvalues of the F-ratio $\mathbf{F} \equiv$ $\mathbf{E}^{-1} \mathbf{H}$ of two large-dimensional Wishart matrices. The experiments are parameterized by the values of a finite number $k$ of spikes that describe the "ratio" of the covariance parameters of $\mathbf{H}$ and $\mathbf{E}$, or, in the case of equal covariance parameters, the non-centrality parameter of $\mathbf{H}$.

We find that the asymptotic behavior of the log ratio of the joint density of the eigenvalues of $\mathbf{F}$, which corresponds to super-critical spikes, to their joint density under a local deviation from these values depends only on the $k$ of the largest eigenvalues $\lambda_{p 1}, \ldots, \lambda_{p k}$. This implies, in particular, that the best statistical inference about $k$ super-critical spikes in the local asymptotic regime is based on the $k$ largest eigenvalue only. A small-scale Monte Carlo analysis shows that $L A N$ confidence sets for super-critical spikes have good coverage properties even for extremely small values of the dimensionality $p$.

As a by-product of our analysis, we establish the joint asymptotic normality of a few of the largest eigenvalues of $\mathbf{F}$ that correspond to the supercritical spikes. We derive explicit formulas for the phase transition threshold, for the almost sure limits of the super-critical eigenvalues, and for the asymptotic variances of their fluctuations around these limits.

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# SUPPLEMATARY MATERIAL FOR "LOCAL ASYMPTOTIC NORMALITY OF THE SPECTRUM OF HIGH-DIMENSIONAL SPIKED F-RATIOS." 

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> This note contains supplementary material for Dharmawansa et al (2016) (DJO in what follows). It is lined up with sections in the main text to make it relatively easy to see how and where the proof details fit in.

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2. Setup. There is no supplementary material for the Setup section of DJO.

## 3. Phase transition and almost sure limits.

3.1. An extension of Bai and Silverstein's (1998) Theorem 1.1. In this subsection, we will use notations from Bai and Silverstein (1998). Their Theorem 1.1 covers only nonrandom $T_{n}$. Remark 6.5 on page 125 of Bai and Silverstein (2010) points out that the theorem is easily extendable to the cases of random $T_{n}$, as long as it is independent from $X_{n}$, its limiting spectral distribution is nonrandom, and condition (f) of the theorem that "interval $[a, b]$ with $a>0$ lies outside the suport
of $F^{c, H}$ and $F^{c_{n}, H_{n}}$ for all large $n$ " holds a.s. (here Bai and Silverstein's $c$ and $c_{n}$ correspond to our $\gamma_{1}$ and $c_{1}$, respectively.)

This still assumes that $\left\|T_{n}\right\|$ is bounded (assumption (d) of the theorem). Unfortunately, if $T_{n}=E^{-2}$, where $E$ is a sample covariance matrix, as in DJO, $\left\|T_{n}\right\|$ can be larger than any positive number with a small probability. Note however, that as long as $\gamma_{2} \in(0,1),\left\|E^{-2}\right\|$ a.s. converges to a finite number. This helps because, as the following lemma shows, Theorem 1.1 remains valid if assumption (d) is replaced by
(d*) There exists $C>0$, such that $\limsup _{p, n \rightarrow \gamma}\left\|T_{n}\right\|<C$, a.s.
Lemma 1. Assumption (d) of Theorem 1.1 can be replaced by ( $d^{*}$ ) without changing the validity of the theorem.

Proof: Define event

$$
\Omega=\left\{\left\|E^{-1}\right\|<C \text { for all } p>p_{0}\right\} .
$$

By ( $\mathrm{d}^{*}$ ), $\operatorname{Pr}(\Omega) \rightarrow 1$ as $p_{0} \rightarrow \infty$. On the other hand, Theorem 1.1 holds conditionally on $\Omega$. That is, the conditional probability $\operatorname{Pr}\left(\Omega_{1} \mid \Omega\right)=1$, where $\Omega_{1}$ is the event that there exists $p_{1}$ s.t. for any $p>p_{1}$, all the eigenvalues of $E^{-1} H$ do not belong to $[a, b]$. Therefore, $\operatorname{Pr}\left(\Omega_{1} \cap \Omega\right)=\operatorname{Pr}(\Omega) \rightarrow 1$ as $p_{0} \rightarrow \infty$. But $\operatorname{Pr}\left(\Omega_{1}\right) \geq \operatorname{Pr}\left(\Omega_{1} \cap \Omega\right)$. Hence, we must have $\operatorname{Pr}\left(\Omega_{1}\right)=1$.

Lemma 1 implies that the largest eigenvalue of $E^{-1} H$ a.s. converges to $\beta_{+}$.
3.2. Proof of Lemma DJO1 about the convergence of $\operatorname{tr}(H-x E)^{-1} / p$. Let $x \in \mathbb{R}$ be such that $x>\beta_{+}$, and let $\hat{F}(\lambda ; x)$ be the empirical distribution function of the eigenvalues of $H-x E$. For any $z \in \mathbb{C}^{+}$, let

$$
\hat{m}(z ; x)=\int(\lambda-z)^{-1} \mathrm{~d} \hat{F}(\lambda ; x)
$$

be the Stieltjes transform of $\hat{F}(\lambda ; x)$. Note that matrix $H-x E$ can be represented in the form $Y T Y^{\prime} / p$, where $Y$ is a $p \times\left(n_{1}+n_{2}\right)$ matrix with i.i.d. standard normal entries and $T$ is a diagonal matrix with the first $n_{1}$ and the last $n_{2}$ diagonal elements equal to $c_{1}$ and $-x c_{2}$, respectively. Therefore, by Theorem 1.1 of Silverstein and Bai (1995), for any $z \in \mathbb{C}^{+}, \hat{m}(z ; x)$ a.s. converges to $m(z ; x) \in \mathbb{C}^{+}$, which is an analytic function in the domain $z \in \mathbb{C}^{+}$that solves the functional equation (DJO13).

By Lemma 1, the largest eigenvalue of $E^{-1} H$ a.s. converges to $\beta_{+}$. Therefore, for any $x>\beta_{+}$, the largest eigenvalue of $H-x E$ is a.s. asymptotically bounded away from the positive semi-axis. Hence, $\hat{m}(z ; x)$ is analytic and bounded in a small disc $D$ around $z=0$ for all sufficiently large $p$ and $n$, a.s. By Vitali's theorem (see Titchmarsh (1939), p. 168), $\hat{m}(z ; x)$ is a.s. converging to an analytic function in $D$. Since, in $D \cap \mathbb{C}^{+}$, the limiting function is $m(z ; x)$, we have

$$
\operatorname{tr}(H-x E)^{-1} / p=\hat{m}(0 ; x) \xrightarrow{\text { a.s. }} m(0 ; x),
$$

where $m(0 ; x)=\lim _{z \rightarrow 0} m(z ; x)$. Further, $\operatorname{tr}(H-\zeta E)^{-1} / p$ is an analytic bounded function of $\zeta$ in a small disk $D_{x}$ around $x$, for all sufficiently large $p$ and $n$, a.s. Therefore, by Vitali's theorem its a.s. limit $f(\zeta)$ is analytic in $D_{x}$, and

$$
\frac{\mathrm{d}}{\mathrm{~d} \zeta} \operatorname{tr}(H-\zeta E)^{-1} / p \xrightarrow{\text { a.s. }} \frac{\mathrm{d}}{\mathrm{~d} \zeta} f(\zeta)
$$

in $D_{x}$. On the other hand, we know that $f(\zeta)=m(0 ; \zeta)$ for $\zeta \in \mathbb{R} \cap D_{x}$. Therefore, we have (DJO12).
3.3. Proof of Lemma DJO2 about the asymptotic proportionality of $M(x)$ and $s+\gamma_{1} I_{k}$. The convergences stated in Lemma DJO2 follow from (DJO5), (DJO9), and Lemma 11 stated below.

Lemma 2. Let $C$ be a random $p \times p$ matrix, independent from $u$ and $v$, which are as defined in Section DJO2, and such that $p\|C\|$ is bounded for all sufficiently large $p$, a.s. Then, as $p \rightarrow \infty$,

$$
\left\|v^{\prime} C v-(\operatorname{tr} C) I_{k}\right\| \xrightarrow{\text { a.s. }} 0 \text { and }\left\|v^{\prime} C u\right\| \xrightarrow{\text { a.s. }} 0 .
$$

Proof: This lemma follows from the Borel-Cantelli lemma, and the upper bounds on the fourth moments of the entries $v^{\prime} C v-(\operatorname{tr} C) I_{k}$ and $v^{\prime} C u$ established by Lemma 2.7 of Bai and Silverstein (1998).
3.4. Proof of Lemma DJO3 about properties of functions $M(x)$ and $m(0 ; x)$. Let $\mu_{1} \in(0, \infty)$ be the largest eigenvalue of $E^{-1} H$. For any $x_{1}>x_{2}>\mu_{1}$, matrix $\left(H-x_{1} E\right)^{-1}-\left(H-x_{2} E\right)^{-1}$ is positive definite, a.s. Part (i) follows from this, from the definition (DJO9) of $M(x)$, and from the fact that $\mu_{1} \xrightarrow{\text { a.s. }} \beta_{+}$. Part (i) together with Lemmas DJO1 and DJO2 imply that $m(0 ; x)$ is increasing on $\left(\beta_{+}, \infty\right)$. It is strictly increasing because, otherwise, equation (DJO13) would not be satisfied for some $z \in \mathbb{C}^{+}$that are sufficiently close to zero. The continuity follows from the analyticity of $m(0 ; x)$ established in the proof of Lemma DJO1. Finally, $\lim _{x \rightarrow \infty} m(0 ; x)=0$ is implied by (ii) and (DJO11). Equation (DJO13) implies that

$$
\lim _{x \downarrow \beta_{+}} m(0 ; x)=\left(\gamma_{2}-1\right) /[(\rho+1) \rho],
$$

which, in its turn, implies the second statement of (iii).

## 4. Asymptotic normality.

4.1. Proof of Lemma DJO5 about $\sigma\left(l_{i}\right)$ and $\delta\left(l_{i}\right)$. By Lemmas DJO1 and DJO2,

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} x} M\left(l_{i}\right) \xrightarrow{\text { a.s. }}\left(s+\gamma_{1} I_{k}\right) \frac{\mathrm{d}}{\mathrm{~d} x} m\left(0 ; x_{i}\right) . \tag{1}
\end{equation*}
$$

Further,

$$
\begin{equation*}
\left(1+M_{i i}\left(l_{i}\right)\right)\left(I_{k}+M\left(l_{i}\right)\right)^{-1} \xrightarrow{\text { a.s. }} \operatorname{diag}\{0, \ldots, 0,1,0, \ldots, 0\} \tag{2}
\end{equation*}
$$

with 1 at the $i$-th place on the diagonal. The latter convergence follows from the fact that $I_{k}+M\left(l_{i}\right)$ can be viewed as a small perturbation of a diagonal matrix

$$
I_{k}+\left(s+\gamma_{1} I_{k}\right) m\left(0 ; x_{i}\right),
$$

which has non-zero diagonal elements, except at the $i$-th position. The eigenvalue perturbation formulae (see, for example, (2.33) on p. 79 of Kato (1980)) will then lead to (2). Combining (10) and (2), and using the definition of $\sigma\left(l_{i}\right)$, we obtain (i).

To establish (ii), we note that $\left(1+M_{i i}\left(l_{i}\right)\right) \operatorname{tr} R\left(l_{i}\right)=O_{\mathrm{P}}(1)$ by an argument similar to that used to establish (i). Further, $\left(\operatorname{tr} J\left(l_{i}\right)\right)^{2}-\operatorname{tr} J^{2}\left(l_{i}\right)$ is a linear function of the only eigenvalue of $J\left(l_{i}\right)$ that diverges to infinity. By the eigenvalue perturbation formulae, such an eigenvalue equals $\left(1+M_{i i}\left(l_{i}\right)\right)^{-1} O(1)$ a.s. Therefore,

$$
\left(1+M_{i i}\left(l_{i}\right)\right)\left(\left(\operatorname{tr} J\left(l_{i}\right)\right)^{2}-\operatorname{tr} J^{2}\left(l_{i}\right)\right)=O(1)
$$

which concludes the proof of (ii).
4.2. Proof of Lemma DJO6 that $Z^{(7)}$ a.s. converges to zero. Recall that

$$
\begin{equation*}
Z^{(7)}=\left(s+c_{1} I_{k}\right) \sqrt{p}\left(\operatorname{tr} G-m_{p}\left(0 ; l_{i}\right)\right) \tag{3}
\end{equation*}
$$

where

$$
G=\left(H-l_{i} E\right)^{-1} / p, l_{i}=\frac{\left(s_{i}+c_{1}\right)\left(s_{i}+1\right)}{s_{i}-c_{2}\left(s_{i}+1\right)}
$$

and $m_{p}\left(z ; l_{i}\right)$ is the Stieltjes transform of the limiting spectral distribution of $H-l_{i} E$ as $n_{1}, n_{2}$, and $p$ diverge to infinity so that $c_{1}$ and $c_{2}$ (and thus $l_{i}$ too) remain fixed.
Let $\hat{F}(\lambda ; x)$ be the empirical distribution of the eigenvalues of $H-x E$ and $\hat{m}(z ; x)$ be its Stieltjes transform. That is,

$$
\hat{m}(z ; x) \equiv \int(\lambda-z)^{-1} \mathrm{~d} \hat{F}(\lambda ; x) \equiv \operatorname{tr}\left(H-x E-z I_{p}\right)^{-1} / p .
$$

Then, to establish Lemma DJO6, it is sufficient to prove that

$$
\begin{equation*}
\sqrt{p}\left(\hat{m}\left(0 ; l_{i}\right)-m_{p}\left(0 ; l_{i}\right)\right) \xrightarrow{\text { a.s. }} 0 . \tag{4}
\end{equation*}
$$

Remark 3. By definition, $l_{i} \rightarrow x_{i}$, where $x_{i}$ is the a.s. limit of the $i$-th largest super-critical eigenvalue of $\mathbf{F}$. Therefore, $l_{i}>\beta_{+}$for sufficiently large $p$, where $\beta_{+}$is the a.s. limit of the largest eigenvalue of $E^{-1} H$ as $p, n \rightarrow_{\gamma} \infty$. This fact implies that, $H-l_{i} E$ is a.s. negative definite for sufficiently large $p$, and $\hat{m}\left(z ; l_{i}\right)$ amd $m_{p}\left(z ; l_{i}\right)$ are well defined for $z=0$ by analytic continuation form the upper half of the complex plane.

Equations for $m_{p}$. By the definition of $H$ and $E$,

$$
H-l_{i} E=Y T Y^{\prime} / p,
$$

where $Y \sim N\left(0, I_{p} \otimes I_{n_{1}+n_{2}}\right)$ and $T=\operatorname{diag}\left(t_{1}, \ldots, t_{n_{1}+n_{2}}\right)$ with

$$
t_{j}=\left\{\begin{array}{c}
c_{1} \text { for } j \leq n_{1} \\
-l_{i} c_{2} \text { for } n_{1}<j \leq n_{1}+n_{2}
\end{array} .\right.
$$

Denote the empirical distribution of $\left\{t_{1}, \ldots, t_{n_{1}+n_{2}}\right\}$ as $\hat{\mathcal{T}}(t)$.
Results of Silverstein and Bai (1995) imply that, as $p, n \rightarrow_{\gamma} \infty, \hat{F}\left(\lambda ; l_{i}\right)$ a.s. weakly converges to a distribution $F^{\gamma, \mathcal{T}}(\lambda)$, whose Stieltjes transform $m(z)$ satisfies equation (compare to equation DJO13)

$$
\begin{equation*}
m(z)=-\left(z-\frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t \mathrm{~d} \mathcal{T}(t)}{1+\operatorname{tm}(z)}\right)^{-1} \tag{5}
\end{equation*}
$$

where $\mathcal{T}(t)$ is the limit of $\hat{\mathcal{T}}(t)$. More explicitly, $m(z)$ satisfies

$$
\begin{equation*}
m(z)=-\left(z-\frac{1}{1+\gamma_{1} m(z)}+\frac{x_{i}}{1-x_{i} \gamma_{2} m(z)}\right)^{-1} . \tag{6}
\end{equation*}
$$

Function $m_{p}(z) \equiv m_{p}\left(z ; l_{i}\right)$ is the analogue of $m(z)$ uner the fixed $p / n_{1} \equiv c_{1}$ and $p / n_{2} \equiv c_{2}$ asymptotics. That is, it satisfies equation

$$
\begin{equation*}
m_{p}(z)=-\left(z-\frac{c_{1}+c_{2}}{c_{1} c_{2}} \int \frac{t \mathrm{~d} \hat{\mathcal{T}}(t)}{1+t m_{p}(z)}\right)^{-1} . \tag{7}
\end{equation*}
$$

For future reference, notice that

$$
\begin{equation*}
z m_{p}(z)-\frac{c_{1}+c_{2}}{c_{1} c_{2}} \int \frac{t m_{p}(z) \mathrm{d} \hat{\mathcal{T}}(t)}{1+t m_{p}(z)}+1=0 \tag{8}
\end{equation*}
$$

REmark 4. Had matrix $T=\operatorname{diag}\left(t_{1}, \ldots, t_{n}\right)$ been positive semi-definite, our 'target' equation (4) would have been following from results of Bai and Silverstein (2004). Our strategy of the proof of (4) will be to extend some of Bai and Silverstein's (2004) analysis to cover $T$ that are not positive semi-definite.

An upper bound on $\left|\sqrt{p}\left(\hat{m}\left(0 ; l_{i}\right)-m_{p}\left(0 ; l_{i}\right)\right)\right|$. Suppose that $\mathcal{K}$ is a contour in the complex plane that does not encircle zero, but encircles all the eigenvalues of $H-l_{i} E$ and the support of the limiting (under fixed $c_{1}, c_{2}$ asymptotics) spectral distribution $F^{c, \hat{\mathcal{T}}}$ of $H-l_{i} E$. Then, we have

$$
\oint_{\mathcal{K}} \frac{\hat{m}\left(z ; l_{i}\right)-m_{p}\left(z ; l_{i}\right)}{z} \mathrm{~d} z=\oint_{\mathcal{K}} \int \frac{\mathrm{d} \hat{F}\left(\lambda ; l_{i}\right)-\mathrm{d} F^{c, \hat{\mathcal{T}}}(\lambda)}{z(\lambda-z)} \mathrm{d} z
$$

Intechanging the order of the integrals, multiplying by $-1 / 2 \pi \mathrm{i}$, and using Cauchy's residue theorem, we obtain

$$
\begin{equation*}
-\frac{1}{2 \pi \mathrm{i}} \oint_{\mathcal{K}} \frac{\hat{m}\left(z ; l_{i}\right)-m_{p}\left(z ; l_{i}\right)}{z} \mathrm{~d} z=\hat{m}\left(0 ; l_{i}\right)-m_{p}\left(0 ; l_{i}\right) \tag{9}
\end{equation*}
$$

To prove the a.s. convergence of $\sqrt{p}\left(\hat{m}\left(0 ; l_{i}\right)-m_{p}\left(0 ; l_{i}\right)\right)$ to zero, we will analyze the behavior of $\sqrt{p}\left(\hat{m}\left(z ; l_{i}\right)-m_{p}\left(z ; l_{i}\right)\right)$ along $\mathcal{K}$. But first, let us explicitly construct such a contour.

Since $l_{i}$ converges to $x_{i}$, which lies above $\beta_{+}$, we will assume without loss of generality that $l_{i} \in\left[x_{-} x_{+}\right]$, where $x_{+}$and $x_{-}$are fixed real numbers satisfying the inequality $x_{+}>x_{-}>\beta_{+}$. Let $a_{-}$be the lower bound of the support of the limiting spectral distribution (LSD) of $H-x_{+} E$ and let $a_{+}$be the upper bound of the support of the LSD of $H-x_{-} E$. Note that, almost surely, $H-x_{+} E<H-x_{-} E<0$ for sufficiently large $p$, and thus, $a_{-} \leq a_{+} \leq 0$. In fact, it easy to see that $a_{+}<0$.

Lemma $5 . a_{+}<0$.
Proof: Let $A(x)=H-x E$. Consider the following decomposition

$$
A(x)=\left(H-\frac{x+\beta_{+}}{2} E\right)-\frac{x-\beta_{+}}{2} E \equiv A_{1}(x)+A_{2}(x)
$$

For the largest eigenvalue of $A_{2}\left(x_{-}\right)$, we have

$$
\begin{equation*}
\lambda_{\max }\left(A_{2}\left(x_{-}\right)\right) \stackrel{a . s .}{\rightarrow}-\frac{x_{-}-\beta_{+}}{2}\left(1-\sqrt{\gamma_{2}}\right)^{2} \equiv \bar{a}_{+}<0 . \tag{10}
\end{equation*}
$$

On the other hand, since $\beta_{+}$is the a.s. limit of the largest eigenvalue of $E^{-1} H, A_{1}\left(x_{-}\right)$is a.s. negative semi-definite for sufficiently large $p$. Hence, $\lambda_{\max }\left(A\left(x_{-}\right)\right)<\bar{a}_{+}$for sufficiently large $p$, a.s., which implies that $a_{+}<\bar{a}_{+}<0$. $\square$

We say that a sequence of events $Q_{p}$ occurs with overwhelming probability (w.o.p.) if $\operatorname{Pr}\left(Q_{p}^{c}\right)=$ $o\left(p^{-t}\right)$ for each fixed $t>0$. Often, we will simply say that $Q_{p}$ occurs w.o.p. omitting the words "the sequence of events".

Lemma 6. The sequence of events

$$
Q_{p}=\left\{\lambda_{\max }\left(A\left(l_{i}\right)\right)<\bar{a}_{+}, \lambda_{\min }\left(A\left(l_{i}\right)\right)>-4 x_{+}\right\}
$$

occurs w.o.p.
Proof: Consider the sequene of events

$$
\hat{Q}_{p}=\left\{\lambda_{\max }\left(A\left(x_{-}\right)\right)<\bar{a}_{+}, \lambda_{\min }\left(A\left(x_{+}\right)\right)>-4 x_{+}\right\} .
$$

Since $l_{i} \in\left[x_{-}, x_{+}\right], \hat{Q}_{p} \subseteq Q_{p}$ and it is sufficient to prove that $\hat{Q}_{p}$ occurs w.o.p.
The decomposition

$$
A(x)=\left(H-\frac{x+3 \beta_{+}}{4} E\right)-\frac{3\left(x-\beta_{+}\right)}{4} E
$$

and the definition $A(x)=H-x E$ show that event $\hat{Q}_{p}^{c}$ implies $\cup_{i=1}^{4} R_{p i}$, where

$$
\begin{aligned}
& R_{p 1}=\left\{\lambda_{\max }\left(-\frac{3\left(x_{-}-\beta_{+}\right)}{4} E\right) \geq \bar{a}_{+}\right\}=\left\{\lambda_{\min }(E) \leq \frac{2}{3}\left(1-\sqrt{\gamma_{2}}\right)^{2}\right\}, \\
& R_{p 2}=\left\{H-\frac{x_{-}+3 \beta_{+}}{4} E<0\right\}^{c}=\left\{\lambda_{\max }\left(E^{-1} H\right) \geq \frac{x_{-}+3 \beta_{+}}{4}\right\}, \\
& R_{p 3}=\left\{\lambda_{\min }\left(-x_{+} E\right) \leq 4 x_{+}\right\}=\left\{\lambda_{\max }(E) \geq 4\right\}, \\
& R_{p 4}=\{H>0\}^{c}=\left\{\lambda_{\min }(H) \leq 0\right\} .
\end{aligned}
$$

The Gaussian concentration inequalities for the largest and smallest singular values of Wishart matrices imply that $R_{p 1}^{c}, R_{p 3}^{c}$, and $R_{p 4}^{c}$ occur w.o.p. Further, as follows, for example, from the proof of Theorem 11.3.2 in Muirhead (1982), the largest root of the equation

$$
\begin{equation*}
\operatorname{det}\left\{n_{1} H-x\left(n_{1} H+n_{2} E\right)\right\}=0 \tag{11}
\end{equation*}
$$

is distributed as the first squared sample canonical correlation coefficient $r_{1}^{2}$ between columns of $Z_{1}$ and $Z_{2}$, where $Z_{1}$ and $Z_{2}$ are independent $p \times\left(n_{1}+n_{2}\right)$ and $n_{1} \times\left(n_{1}+n_{2}\right)$ matrices with independent $N(0,1)$ entries. In the next subsection of this note, we show that such a squared sample canonical correlation coefficient satisfies the following concentration inequality

$$
\begin{equation*}
\operatorname{Pr}\left\{r_{1}^{2}>\mathbb{E} r_{1}^{2}+t\right\} \leq 2 \exp \left\{-\frac{\left(n_{1}+n_{2}\right) t^{2}}{2 \times 16^{2}}\right\}, t>0 \tag{12}
\end{equation*}
$$

This probability bound is not the best possible, but sufficient for our purposes. Indeed, note that the largest root of (11) equals

$$
\lambda_{\max }\left(E^{-1} H\right) /\left(n_{2} / n_{1}+\lambda_{\max }\left(E^{-1} H\right)\right)
$$

This equality, the fact that $\lambda_{\max }\left(E^{-1} H\right)$ a.s. converges to $\beta_{+}$, and the concentration inequality (12) imply that $R_{p 2}^{c}$ occurs w.o.p. Since $\hat{Q}_{p}$ is implied by $\cap_{i=1}^{4} R_{p i}^{c}, \hat{Q}_{p}$ also occurs w.o.p. $\square$

Remark 7. The bounds $-4 x_{+}$and $\bar{a}_{+}$on the smallest and the largest eigenvalues of $A\left(l_{i}\right)$ are rough, but they are sufficient for our purposes.

Now we are ready to construct contour $\mathcal{K}$. It is the rectangle shown in Figure 1. The contour intersects the real axis at $-5 x_{+}$and $\bar{a}_{+} / 2$, so that the bounds $-4 x_{+}$and $\bar{a}_{+}$remain inside the contour, but zero lies outside the contour. It is symmetric around the real axis.

Separate, but related, arguments for bounding $\sqrt{p}\left(\hat{m}\left(z ; l_{i}\right)-m_{p}\left(z ; l_{i}\right)\right)$ are needed for the horizontal and vertical segments of the contour, $\mathcal{K}_{H}$ and $\mathcal{K}_{V}$ respectively. Small vertical intervals $\mathcal{K}_{0 p}=\left\{z \in \mathcal{K}:|\operatorname{Im} z| \leq p^{-2}\right\}$ about the real axis will be excluded from many bounds and handled separately. Accordingly, we write $\mathcal{K}_{p}$ for $\mathcal{K} \backslash \mathcal{K}_{0 p}$ and $\mathcal{K}_{V, p}$ for $\mathcal{K}_{V} \backslash \mathcal{K}_{0 p}$. Without loss of generality, we set $|\operatorname{Im} z|=\min \left\{-\bar{a}_{+} / 2, x_{+}\right\} \equiv \eta$ for $z \in \mathcal{K}_{H}$. The purpose of such a setting is to have a distance between $\left[-4 x_{+}, \bar{a}_{+}\right]$and $\gamma$ be bounded from below by $\eta$.


Fig 1. Contour $\mathcal{K}$ in $u+i v$ plane
Define 'deterministic' and 'stochastic' terms by

$$
\begin{aligned}
I_{p}^{(D)} & =\oint_{\mathcal{K}_{p}} \sqrt{p}\left|\mathbb{E} \hat{m}\left(z ; l_{i}\right)-m_{p}\left(z ; l_{i}\right)\right| \mathrm{d} z \text { and } \\
I_{p}^{(S)} & =\oint_{\mathcal{K}_{p}} \sqrt{p}\left|\hat{m}\left(z ; l_{i}\right)-\mathbb{E} \hat{m}\left(z ; l_{i}\right)\right| \mathrm{d} z
\end{aligned}
$$

and an exceptional term near the real axis by

$$
I_{p}^{(E)}=\oint_{\mathcal{K}_{0 p}} \sqrt{p}\left|\hat{m}\left(z ; l_{i}\right)-m_{p}\left(z ; l_{i}\right)\right| \mathrm{d} z
$$

Write $\left\|z^{-1}\right\|_{\infty, \mathcal{K}}=\sup \left\{z^{-1}: z \in \mathcal{K}\right\}$. From (9), we have on event $Q_{p}$,

$$
\begin{equation*}
\left|\sqrt{p}\left(\hat{m}\left(0 ; l_{i}\right)-m_{p}\left(0 ; l_{i}\right)\right)\right| \leq(2 \pi)^{-1}\left\|z^{-1}\right\|_{\infty, \mathcal{K}}\left\{I_{p}^{(D)}+I_{p}^{(S)}+I_{p}^{(E)}\right\} \tag{13}
\end{equation*}
$$

First reduction. Let us show that the proof of the convergence (4) can be reduced to verifying the stochastic bounds

$$
\begin{align*}
& \sup _{z \in \mathcal{K}_{p}} \mathbb{E}\left|\hat{m}\left(z ; l_{i}\right)-\mathbb{E} \hat{m}\left(z ; l_{i}\right)\right|^{2 l} \leq C_{l} p^{-2 l}, l=1,2  \tag{14}\\
& \text { imsart-aos ver. 2012/04/10 file: supercriticalsm1.tex date: January 28, } 2017
\end{align*}
$$

and the deterministic convergence

$$
\begin{equation*}
\sup _{z \in \mathcal{K}_{p}} \sqrt{p}\left|\mathbb{E} \hat{m}\left(z ; l_{i}\right)-m_{p}\left(z ; l_{i}\right)\right| \rightarrow 0 \tag{15}
\end{equation*}
$$

Note that (15) implies that $I_{p}^{(D)} \rightarrow 0$.
For the stochastic bounds, write $\left|\mathcal{K}_{p}\right|$ for the length of $\mathcal{K}_{p}$, and make use of Hölder's inequality for $\left(I_{p}^{(S)}\right)^{4}$ and then (14) to bound

$$
\begin{aligned}
\operatorname{Pr}\left(I_{p}^{(S)}>\varepsilon\right) & \leq \varepsilon^{-4} \mathbb{E}\left(I_{p}^{(S)}\right)^{4} \\
& \leq \varepsilon^{-4}\left|\mathcal{K}_{p}\right|^{3} \oint_{\mathcal{K}_{p}} p^{2} \mathbb{E}\left|\hat{m}\left(z ; l_{i}\right)-\mathbb{E} \hat{m}\left(z ; l_{i}\right)\right|^{4}|\mathrm{~d} z| \\
& \leq \varepsilon^{-4}\left|\mathcal{K}_{p}\right|^{4} p^{2} \sup _{z \in \mathcal{K}_{p}} \mathbb{E}\left|\hat{m}\left(z ; l_{i}\right)-\mathbb{E} \hat{m}\left(z ; l_{i}\right)\right|^{4} \\
& \leq \varepsilon^{-4}\left|\mathcal{K}_{p}\right|^{4} C_{2} p^{-2} .
\end{aligned}
$$

Since this sequence is summable in $p$, we have $I_{p}^{(S)} \rightarrow 0$, a.s.
We turn to the exceptional term $I_{p}^{(E)}$. When $z \in \mathcal{K}_{0 p}$ and event $Q_{p}$ occurs, we may bound $\hat{m}\left(z ; l_{i}\right)=\frac{1}{p} \sum_{j=1}^{p}\left(\mu_{p j}-z\right)^{-1}$ (where $\mu_{p j}$ with $j=1, \ldots, p$ are the eigenvalues of $A\left(l_{i}\right) \equiv H-l_{i} E$ ) using

$$
\max _{j}\left|\mu_{p j}-z\right|^{-1} \leq \max \left\{-2 / \bar{a}_{+}, 1 / x_{+}\right\} \equiv \eta^{-1}
$$

so that $\left|\hat{m}\left(z ; l_{i}\right)\right| \leq \eta^{-1}$. Further, for sufficiently large $p$, we have $\left|m_{p}\left(z ; l_{i}\right)\right| \leq \eta^{-1}$ for $z \in \mathcal{K}_{0 p}$. Consequently,

$$
\begin{aligned}
\operatorname{Pr}\left\{I_{p}^{(E)}>\varepsilon, Q_{p}\right\} & \leq \varepsilon^{-1} \mathbb{E}\left(I_{p}^{(E)} \mathbf{1}_{Q_{p}}\right) \\
& \leq \varepsilon^{-1}\left|\mathcal{K}_{0_{p}}\right| \sup _{z \in \mathcal{K}_{0_{p}}}\left\{\sqrt{p}\left(\left|\hat{m}\left(z ; l_{i}\right)\right|+\left|m_{p}\left(z ; l_{i}\right)\right|\right) \mathbf{1}_{Q_{p}}\right\} \\
& \leq \varepsilon^{-1} 4 p^{-2} 2 \eta^{-1} p^{1 / 2}=C \varepsilon^{-1} p^{-3 / 2}
\end{aligned}
$$

where $C$ denotes a constant. Again this is summable in $p$, and since $Q_{p}$ occurs w.o.p., it follows that $I_{p}^{(E)} \xrightarrow{\text { a.s. }} 0$ also. In summary, referring to (13), we see that in order to show a.s. convergence in (4), it remains to establish (14) and (15).

We begin with some preliminary results. Two tools for handling fluctuations are then introduced: first, moment bounds for deviations of quadratic forms, and then, the martingale difference structure. Then we proceed to bound the deterministic term in (13). After all this, we are ready to attack the stochastic bounds (14).

Preliminary results. The approach consists in careful analysis of the perturbations induced in the resolvent of $A\left(l_{i}\right)=Y T Y^{\prime} / p$ by deletion of a single column from $Y$. Thus, let $q_{k}=Y_{\cdot k} / \sqrt{p}$ (the $k$-th column of $Y$ divided by $\sqrt{p}$ ) and $A_{k}=A\left(l_{i}\right)-t_{k} q_{k} q_{k}^{\prime}$. Consider events

$$
Q_{p k}=\left\{\lambda_{\max }\left(A_{k}\right)<\bar{a}_{+}, \lambda_{\min }\left(A_{k}\right)>-4 x_{+}\right\} .
$$

Similar to $Q_{p}$, events $Q_{p k}$ occur w.o.p.
Let

$$
Q_{p, a l l}=Q_{p} \cap Q_{p 1} \cap \ldots \cap Q_{p n} .
$$

Then, $Q_{p, \text { all }}$ occur w.o.p. This follows from the equality $\max _{j=1, \ldots, p} \operatorname{Pr}\left\{Q_{p j}^{c}\right\}=o\left(p^{-t}\right)$ for each fixed $t>0$. The equality is true because, first, each $Q_{p j}, j=1, \ldots, p$ occurs w.o.p., and, second, $\operatorname{Pr}\left\{Q_{p j}^{c}\right\}$ takes on only two possible values, depending on whether $t_{k}=c_{1}$ or $t_{k}=-l_{i} c_{2}$. In particular, although there is a proliferation of the number of events involved in the construction of $Q_{p, a l l}$ as $p \rightarrow \infty$, the probabilities of these events approach one uniformly.

Let

$$
D=\left(A\left(l_{i}\right)-z I_{p}\right)^{-1} \text { and } D_{k}=\left(A_{k}-z I_{p}\right)^{-1},
$$

and let

$$
\beta_{k}=1+t_{k} q_{k}^{\prime}\left(A_{k}-z I\right)^{-1} q_{k} .
$$

An important identity to be used later is

$$
\begin{equation*}
q_{k}^{\prime} D=q_{k}^{\prime} D_{k} / \beta_{k} . \tag{16}
\end{equation*}
$$

The following lemma establishes a useful bound on $\beta_{k}$.
Lemma 8. Suppose that event $Q_{p, \text { all }}$ holds. Then, there exists a constant $c$, that depends only on $\bar{a}_{+}, x_{+}$, and $\mathcal{K}$, such that, for any $z=u+i v \in \mathcal{K}$,

$$
\left|\beta_{k}\right| \geq \sqrt{\frac{v^{2}}{c^{2}+v^{2}}}
$$

Proof: Let $e_{j}$ be a normalized eigenvector corresponding to the $j$-th largest eigenvalue, $\mu_{j k}$, of $A_{k}$. Then, we have

$$
\begin{equation*}
\beta_{k}=1+t_{k} q_{k}^{\prime}\left(A_{k}-z I\right)^{-1} q_{k}=1+t_{k} \sum_{j=1}^{p} \frac{\left(q_{k}^{\prime} e_{j}\right)^{2}}{\mu_{j k}-z} . \tag{17}
\end{equation*}
$$

Consider the case where $v>0$. When $Q_{p, \text { all }}$ holds, $\mu_{j k} \in\left[-4 x_{+}, \bar{a}_{+}\right]$, and

$$
\arg \left(\mu_{j k}-z\right) \in\left[\arg \left(-4 x_{+}-z\right), \arg \left(\bar{a}_{+}-z\right)\right],
$$

where arg belongs to $(-\pi, \pi)$, and $\arg \left(-4 x_{+}-z\right)<\arg \left(\bar{a}_{+}-z\right)<0$. Let us denote $\arg \left(-4 x_{+}-z\right)$ as $-\varphi_{L}$ and $\arg \left(\bar{a}_{+}-z\right)$ as $-\varphi_{R}$. Note that

$$
\arg \left\{\left(\mu_{j k}-z\right)^{-1}\right\} \in\left[\varphi_{R}, \varphi_{L}\right]
$$

and thus,

$$
\varphi_{R} \leq \arg \sum_{j=1}^{p}\left(q_{k}^{\prime} e_{j}\right)^{2} /\left(\mu_{j k}-z\right) \leq \varphi_{L} .
$$

These inequalities and equation (17) imply that, when $Q_{p, \text { all }}$ holds, $\left|\beta_{k}\right|$ cannot be smaller than the distance from the origin to the cone $\left\{1+\rho e^{i \varphi}: \rho \in \mathbb{R}, \varphi_{R} \leq \varphi \leq \varphi_{L}\right\}$, which equals $\min \left\{\sin \varphi_{L}, \sin \varphi_{R}\right\}$. On the other hand,

$$
\sin \varphi_{L}=\sin \arctan \frac{v}{\left|-4 x_{+}-u\right|} \text { and } \sin \varphi_{R}=\sin \arctan \frac{v}{\left|\bar{a} \bar{a}_{+}-u\right|},
$$

so that there exists $c$ that depends only on $\bar{a}_{+}, x_{+}$, and $\mathcal{K}$, such that, when $Q_{p, \text { all }}$ holds,

$$
\left|\beta_{k}\right| \geq \sin \arctan \frac{v}{c}=\sqrt{\frac{v^{2}}{c^{2}+v^{2}}} .
$$

The case where $v<0$ leads to the same conclusion in a similar way.

Remark 9. In the case where all $t_{k}$ are non-negative (which was studied by Bai and Silverstein (2004)), $\left|\beta_{k}\right|$ is always bounded by $|v / z|$. In contrast, when some $t_{k}$ are negative, the corresponding $\left|\beta_{k}\right|$ can be arbitrarily close to zero with small, but positive, probability. This is why the bound in Lemma 8 is conditioned on the event $Q_{p, a l l}$.

For the analysis of the deterministic term $I_{p}^{(D)}$, we will need the following preliminary convergence result.

Lemma 10. $\mathbb{E} \hat{m}\left(z ; l_{i}\right)-m(z) \rightarrow 0$ and $m_{p}\left(z ; l_{i}\right)-m(z) \rightarrow 0$, uniformly in $z \in \mathcal{K}_{p}$. Here $m(z)$ is the Stieltjes transform of the LSD of $A\left(l_{i}\right) \equiv H-l_{i} E$ as $p, n \rightarrow_{\gamma} \infty$.

Proof: For $m_{p}\left(z ; l_{i}\right)$, such a convergence is a consequence of the fact that $F^{c, \hat{\mathcal{T}}}$, that is, the LSD of $A\left(l_{i}\right)$ under the fixed $c_{1}, c_{2}$ asymptotic regime converges to $F^{\gamma, \mathcal{T}}$ as $c_{1} \rightarrow \gamma_{1}$ and $c_{2} \rightarrow \gamma_{2}$. Moreover, the supports of $F^{c, \hat{\mathcal{T}}}$ and $F^{\gamma, \mathcal{T}}$ coinside asymptotically, and lie at a positive distance from contour $\mathcal{K}$.

For $\mathbb{E} \hat{m}\left(z ; l_{i}\right)$, note that, since the spectral distribution $\hat{F}_{p}$ of $A\left(l_{i}\right)$ a.s. converges to $F^{c, \mathcal{T}}$ (we denote this as $\hat{F}_{p} \xrightarrow{d} F^{c, \mathcal{T}}$, a.s.), we have, by the dominated convergence theorem,

$$
\begin{equation*}
\mathbb{E} \hat{F}_{p} \xrightarrow{d} F^{c, \mathcal{T}} . \tag{18}
\end{equation*}
$$

Further, since $\left\{(\lambda-z)^{-1}: z \in \mathcal{K}_{H}\right\}$ is a family of bounded equicontinuous functions of $\lambda \in \mathbb{R}$, (18) implies that $\mathbb{E} \hat{m} \rightarrow m$, uniformly in $z \in \mathcal{K}_{H}$. Next,

$$
\begin{aligned}
\mathbb{E} \hat{m}-m= & \int(\lambda-z)^{-1} \mathbf{1}_{\left[-4 x_{+}, \bar{a}_{+}\right]}(\lambda) \mathrm{d}\left(\mathbb{E} \hat{F}_{p}-F^{c, \mathcal{T}}\right) \\
& +\int(\lambda-z)^{-1} \mathbf{1}_{\left[-4 x_{+}, \bar{a}_{+}\right]^{c}}(\lambda) \mathrm{d} \mathbb{E} \hat{F}_{p}
\end{aligned}
$$

The first integral converges to zero uniformly in $z \in \mathcal{K}_{p V}$ because

$$
\left\{(\lambda-z)^{-1} \mathbf{1}_{\left[-4 x_{+}, \bar{a}_{+}\right]}(\lambda): z \in \mathcal{K}_{p V}\right\}
$$

is a family of bounded equicontinuous functions of $\lambda \in \mathbb{R}$. For the second integral, we have

$$
\sup _{z \in \mathcal{K}_{p V}}\left|\int(\lambda-z)^{-1} \mathbf{1}_{\left[-4 x_{+}, \bar{a}_{+}\right]^{c}}(\lambda) d \mathbb{E} F_{p}\right| \leq p^{2} \operatorname{Pr}\left\{Q_{p}^{c}\right\} \rightarrow 0 .
$$

## Deviations of quadratic forms and the martingale difference structure.

Lemma 11. Let $C$ be a $p \times p$ non-random matrix, and $r \geq 1$. Suppose that $\xi_{j}, j=1, \ldots, p$ are independent mean zero random variables with $\mathbb{E}\left|\xi_{j}\right|^{2}=1$ and $\mathbb{E}\left|\xi_{j}\right|^{l}=\nu_{l}$ for $l \leq 2 r$. Then

$$
\mathbb{E}\left|\xi^{*} C \xi-\operatorname{tr} C\right|^{r} \leq K_{r}\left(\left[\nu_{4} \operatorname{tr}\left(C C^{*}\right)\right]^{r / 2}+\nu_{2 r} \operatorname{tr}\left(C C^{*}\right)^{r / 2}\right)
$$

This is Lemma 2.7 in Bai and Silverstein (1998).
Introduce increasing $\sigma$-fields $\mathcal{F}_{k}=\sigma\left\{q_{j}: j \leq k\right\}$ and let $\mathbb{E}_{k}$ denote $\mathbb{E}\left(\cdot \mid \mathcal{F}_{k}\right)$. For $y=\operatorname{tr} D$, write

$$
y-\mathbb{E} y=\sum_{k=1}^{n_{1}+n_{2}} \mathbb{E}_{k} y-\mathbb{E}_{k-1} y
$$

Now introduce $y_{k}=\operatorname{tr} D_{k}$ and observe that $\mathbb{E}_{k} y_{k}=\mathbb{E}_{k-1} y_{k}$. Therefore,

$$
\begin{gather*}
p\left(\hat{m}\left(z ; l_{i}\right)-\mathbb{E} \hat{m}\left(z ; l_{i}\right)\right)=y-\mathbb{E} y=\sum_{k=1}^{n_{1}+n_{2}} z_{k}, \\
z_{k}=\left(\mathbb{E}_{k}-\mathbb{E}_{k-1}\right)\left(y-y_{k}\right) . \tag{19}
\end{gather*}
$$

We have $\mathbb{E}_{k-1} z_{k}=0$ for $k=1, \ldots, n_{1}+n_{2}$ and so the $z_{k}$ form a martingale difference sequence. The differences are orthogonal, so

$$
\mathbb{E}(y-\mathbb{E} y)^{2}=\sum_{k=1}^{n_{1}+n_{2}} \mathbb{E} z_{k}^{2} .
$$

Let us establish some bounds on $y-y_{k}$. We have

$$
\begin{align*}
D & =\left(A_{k}-z I_{p}+t_{k} q_{k} q_{k}^{\prime}\right)^{-1} \\
& =D_{k}-t_{k} D_{k} q_{k}\left(1+t_{k} q_{k}^{\prime} D_{k} q_{k}\right)^{-1} q_{k}^{\prime} D_{k} \tag{20}
\end{align*}
$$

Therefore,

$$
\begin{equation*}
y-y_{k}=\operatorname{tr}\left(D-D_{k}\right)=-\frac{t_{k}}{1+t_{k} q_{k}^{\prime} D_{k} q_{k}} q_{k}^{\prime} D_{k}^{2} q_{k} . \tag{21}
\end{equation*}
$$

Let $e_{j k}$ be a normalized eigenvector of $A_{k}$ corresponding to its $j$-th largest eigenvalue $\mu_{j k}$. We have

$$
\left|1+t_{k} q_{k}^{\prime} D_{k} q_{k}\right| \geq\left|t_{k}\right|\left|\operatorname{Im} q_{k}^{\prime} D_{k} q_{k}\right|=\left|t_{k}\right| \sum_{j=1}^{p} \frac{\left(q_{k}^{\prime} e_{j k}\right)^{2}}{\left|\mu_{j k}-z\right|^{2}}|v|
$$

and

$$
\left|t_{k} q_{k}^{\prime} D_{k}^{2} q_{k}\right| \leq\left|t_{k}\right| \sum_{j=1}^{p} \frac{\left(q_{k}^{\prime} e_{j k}\right)^{2}}{\left|\mu_{j k}-z\right|^{2}}
$$

Therefore,

$$
\begin{equation*}
\left|y-y_{k}\right| \leq 1 /|v| . \tag{22}
\end{equation*}
$$

This bound can be very large when $v$ is small. Therefore, for $z \in \mathcal{K}_{V}$, we will need another bound. Denote the eigenvalues of $A\left(l_{i}\right)$ as $\mu_{1} \geq \mu_{2} \geq \ldots \geq \mu_{p}$. We would like to show that if $Q_{p, \text { all }}$ holds and $z \in \mathcal{K}_{V}$, we have

$$
\left|y-y_{k}\right|<C
$$

for some constant $C$. There are two vertical sections of $\mathcal{K}_{V}$. Let us denote the "left" vertical section $\mathcal{K}_{V L}$. Consider $z \in \mathcal{K}_{V L}$, let $z_{L}=-5 x_{+}$denote the intercection of $\mathcal{K}_{V L}$ with real axis.

Recall that the eigenvalues of $A_{k}$ are denoted as $\mu_{1 k} \geq \mu_{2 k} \geq \ldots \geq \mu_{p k}$. By interlacing inequality, if $t_{k}>0$, then

$$
\begin{equation*}
\mu_{1} \geq \mu_{1 k} \geq \ldots \geq \mu_{p} \geq \mu_{p k} . \tag{23}
\end{equation*}
$$

If $t_{k}<0$, then

$$
\begin{equation*}
\mu_{1 k} \geq \mu_{1} \geq \ldots \geq \mu_{p k} \geq \mu_{p} \tag{24}
\end{equation*}
$$

Now,

$$
y-y_{k}=\operatorname{Re}\left[\operatorname{tr} D-\operatorname{tr} D_{k}\right]+\mathrm{i} \operatorname{Im}\left[\operatorname{tr} D-\operatorname{tr} D_{k}\right] \equiv u_{1 k}+\mathrm{i} u_{2 k} .
$$

For $u_{1 k}$ and $z \in \mathcal{K}_{V L}$, we have

$$
u_{1 k}=\sum_{j=1}^{p} \frac{\mu_{j}-z_{L}}{\left(\mu_{j}-z_{L}\right)^{2}+v^{2}}-\sum_{s=1}^{p} \frac{\mu_{s k}-z_{L}}{\left(\mu_{s k}-z_{L}\right)^{2}+v^{2}} .
$$

Since $|v| \leq \eta \equiv \min \left\{-\bar{a}_{+} / 2, x_{+}\right\}$, the ratios in the above displayed expression are strictly decreasing functions of $\mu_{j}, \mu_{s k} \geq-4 x_{+}$. Therefore, by interlacing inequalities (23, 24), if $t_{k}>0$, then

$$
\begin{equation*}
-\frac{\mu_{p k}-z_{L}}{\left(\mu_{p k}-z_{L}\right)^{2}+v^{2}}<u_{1 k}<0 . \tag{25}
\end{equation*}
$$

If $t_{k}<0$, then

$$
\begin{equation*}
0<u_{1 k}<\frac{\mu_{p}-z_{L}}{\left(\mu_{p}-z_{L}\right)^{2}+v^{2}} . \tag{26}
\end{equation*}
$$

Similarly,

$$
u_{2 k}=\sum_{j=1}^{p} \frac{v}{\left(\mu_{j}-z_{L}\right)^{2}+v^{2}}-\sum_{s=1}^{p} \frac{v}{\left(\mu_{s k}-z_{L}\right)^{2}+v^{2}}
$$

and the ratios in the above displayed expression are strictly decreasing (increasing) functions of $\mu_{j}, \mu_{s k} \geq-4 x_{+}$when $v>0(v<0)$. Therefore, by (23,24), we have, if $t_{k}>0$,

$$
\begin{equation*}
-\frac{v}{\left(\mu_{p k}-z_{L}\right)^{2}+v^{2}}<u_{2 k}<0 . \tag{27}
\end{equation*}
$$

If $t_{k}<0$, then

$$
\begin{equation*}
0<u_{2 k}<\frac{v}{\left(\mu_{p}-z_{L}\right)^{2}+v^{2}} . \tag{28}
\end{equation*}
$$

From (25-28), we see that, if $Q_{p, \text { all }}$ holds, then, for $z \in \mathcal{K}_{V L}$ with $v \leq \eta$,

$$
\begin{align*}
\left|y-y_{k}\right| & \leq \frac{|v|}{\eta^{2}+v^{2}}+\frac{\eta}{\eta^{2}+v^{2}}  \tag{29}\\
& \leq \frac{\eta}{2 \eta^{2}}+\frac{\eta}{\eta^{2}}=\frac{3}{2} \eta^{-1} .
\end{align*}
$$

Similarly, we can show that the same inequality holds for $z \in \mathcal{K}_{V R}$ (the "right" portion of $\mathcal{K}_{V}$ ).
Another bound on $\beta_{k}$. The bound on $\beta_{k}$ obtained in Lemma 8 will be sufficient for our analysis in cases where $z \in \mathcal{K}_{H}$. However, for $z \in \mathcal{K}_{V}$, it may be too close to zero, and we need another bound.

As follows from (6), for any $t$ from the support of $\mathcal{T}$ (in our case there are only two such $t$ : $\gamma_{1}$ and $\left.-x_{i} \gamma_{2}\right), 1+\operatorname{tm}(z) \neq 0$ for $z \in \mathcal{K}$. Therefore,

$$
\nu \equiv \min _{t \in\left\{\gamma_{1},-x_{i} \gamma_{2}\right\}} \inf _{z \in \mathcal{K}}|1+\operatorname{tm}(z)|>0 .
$$

Since $\mathbb{E} \hat{m}\left(z ; l_{i}\right)-m(z) \rightarrow 0$ uniformly in $z \in \mathcal{K}_{p}$, and since $c_{1}, c_{2}, l_{i} \rightarrow \gamma_{1}, \gamma_{2}, x_{i}$, we have

$$
\min _{t \in\left\{c_{1},-l_{i} c_{2}\right\}} \inf _{z \in \mathcal{K}_{p}}\left|1+t \mathbb{E} \hat{m}\left(z ; l_{i}\right)\right|>2 \nu / 3
$$

for sufficiently large $p$.
Consider the event $\left|\beta_{k}\right| \leq \nu / 3$. If this event holds, then $\left|\beta_{k}-1-t_{k} \mathbb{E} \hat{m}\left(z ; l_{i}\right)\right|>\nu / 3$ for sufficiently large $p$ and any $z \in \mathcal{K}_{p}$. Recalling the definition of $\beta_{k}$, we obtain

$$
\left|t_{k}\right|\left|q_{k}^{\prime} D_{k} q_{k}-\mathbb{E} \hat{m}\left(z ; l_{i}\right)\right| \equiv\left|t_{k}\right|\left|\varepsilon_{k}\right|>\nu / 3
$$

Let us show that the sequence of events $\left|t_{k}\right|\left|\varepsilon_{k}\right| \leq \nu / 3$ occurs w.o.p. Note that this would imply that the sequence of events $\left|\beta_{k}\right|>\nu / 3$ also occurs w.o.p.

We have

$$
\left|\varepsilon_{k}\right| \leq\left|q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right|+\left|\frac{1}{p}\left(\operatorname{tr} D_{k}-\operatorname{tr} D\right)\right|+\left|\frac{1}{p}(\operatorname{tr} D-\mathbb{E} \operatorname{tr} D)\right|
$$

Bound on $\frac{1}{p}(\operatorname{tr} D-\mathbb{E} \operatorname{tr} D)$.
By the Burkholder inequality (see Burkholder, 1973, and Theorem 2.10 of Hall and Heyde, 1980), we have, for any $1<r<\infty$,

$$
\begin{aligned}
& \mathbb{E}|\operatorname{tr} D-\mathbb{E} \operatorname{tr} D|^{r}=\mathbb{E}|y-\mathbb{E} y|^{r} \\
\leq & C_{r} \mathbb{E}\left|\sum_{k=1}^{n_{1}+n_{2}}\left[\left(\mathbb{E}_{k}-\mathbb{E}_{k-1}\right)\left(y-y_{k}\right)\right]^{2}\right|^{r / 2} .
\end{aligned}
$$

Therefore,

$$
\begin{aligned}
& \operatorname{Pr}\left(\frac{\left|t_{k}\right|}{p}|\operatorname{tr} D-\mathbb{E} \operatorname{tr} D|>\frac{\nu}{9}\right) \leq\left(\frac{9\left|t_{k}\right|}{\nu p}\right)^{r} \mathbb{E}|\operatorname{tr} D-\mathbb{E} \operatorname{tr} D|^{r} \\
\leq & \left(\frac{9\left|t_{k}\right|}{\nu p}\right)^{r} C_{r} \mathbb{E}\left|\sum_{k=1}^{n_{1}+n_{2}}\left[\left(\mathbb{E}_{k}-\mathbb{E}_{k-1}\right)\left(y-y_{k}\right)\right]^{2}\right|^{r / 2} .
\end{aligned}
$$

For $z \in \mathcal{K}_{H}$, using (22), we obtain

$$
\operatorname{Pr}\left(\frac{\left|t_{k}\right|}{p}|\operatorname{tr} D-\mathbb{E} \operatorname{tr} D|>\frac{\nu}{9}\right) \leq\left(\frac{9\left|t_{k}\right|}{\nu}\right)^{r} p^{-r / 2} C_{r}\left(\frac{2}{\eta}\right)^{r}\left(\frac{n_{1}+n_{2}}{p}\right)^{r / 2}
$$

Since $r$ is an arbitrary number larger than one, the sequence of events $\frac{\left|t_{k}\right|}{p}|\operatorname{tr} D-\mathbb{E} \operatorname{tr} D| \leq \frac{\nu}{9}$ $k=1,2, \ldots$ occurs w.o.p.

For $z \in \mathcal{K}_{p V}$, we need another estimate. We have

$$
\begin{aligned}
& \operatorname{Pr}\left(\frac{\left|t_{k}\right|}{p}|\operatorname{tr} D-\mathbb{E} \operatorname{tr} D|>\frac{\nu}{9}, Q_{p, a l l}\right) \\
\leq & \operatorname{Pr}\left(\frac{\left|t_{k}\right|}{p}\left|\operatorname{tr} D \mathbf{1}_{Q_{p, a l l}}-\mathbb{E}\left[\operatorname{tr} D\left(\mathbf{1}_{Q_{p, a l l}}+\mathbf{1}_{Q_{p, a l}^{c}}\right)\right]\right|>\frac{\nu}{9}\right)
\end{aligned}
$$

But, for $z \in \mathcal{K}_{p V},|\operatorname{tr} D| \leq p^{3}$, and therefore, $\mathbb{E}\left[\operatorname{tr} D \mathbf{1}_{Q_{p, a l l}^{c}}\right] \leq p^{3} \operatorname{Pr}\left(Q_{p, \text { all }}^{c}\right) \rightarrow 0$ as $p \rightarrow \infty$. In particular, for sufficiently large $p, \frac{\left|t_{k}\right|}{p}\left|\mathbb{E}\left[\operatorname{tr} D \mathbf{1}_{Q_{p, a l l}^{c}}\right]\right|<\frac{\nu}{9}-\frac{\nu}{10}$, and we can write, for sufficiently
large $p$ and any $r>1$,

$$
\begin{aligned}
& \operatorname{Pr}\left(\frac{\left|t_{k}\right|}{p}|\operatorname{tr} D-\mathbb{E} \operatorname{tr} D|>\frac{\nu}{9}, Q_{p, a l l}\right) \\
\leq & \operatorname{Pr}\left(\frac{\left|t_{k}\right|}{p}\left|\operatorname{tr} D \mathbf{1}_{Q_{p, a l l}}-\mathbb{E}\left[\operatorname{tr} D \mathbf{1}_{Q_{p, a l l}}\right]\right|>\frac{\nu}{10}\right) \\
\leq & \left(\frac{10\left|t_{k}\right|}{\nu p}\right)^{r} \mathbb{E}\left|\operatorname{tr} D \mathbf{1}_{Q_{p, a l l}}-\mathbb{E}\left[\operatorname{tr} D \mathbf{1}_{Q_{p, a l l}}\right]\right|^{r} \\
\leq & \left(\frac{10\left|t_{k}\right|}{\nu p}\right)^{r} C_{r} \mathbb{E}\left|\sum_{k=1}^{n_{1}+n_{2}}\left[\left(\mathbb{E}_{k}-\mathbb{E}_{k-1}\right)\left(y \mathbf{1}_{Q_{p, a l l}}-y_{k} \mathbf{1}_{Q_{p, a l l}}\right)\right]^{2}\right|^{r / 2} \\
\leq & \left(\frac{10\left|t_{k}\right|}{\nu}\right)^{r} p^{-r / 2} C_{r}\left(\frac{n_{1}+n_{2}}{p}\right)^{r / 2}\left(3 \eta^{-1}\right)^{r} .
\end{aligned}
$$

The last inequality follows from (29). Since $Q_{p, \text { all }}$ occur w.o.p., the obtained upper bound on $\operatorname{Pr}\left(\frac{\left|t_{k}\right|}{p}|\operatorname{tr} D-\mathbb{E} \operatorname{tr} D|>\frac{\nu}{9}, Q_{p, \text { all }}\right)$ implies that $\frac{\left|t_{k}\right|}{p}|\operatorname{tr} D-\mathbb{E} \operatorname{tr} D| \leq \frac{\nu}{9}$ occur w.o.p. too.

Bound on $q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}$.
By Chebyshev's inequality, for any $r \geq 1$,

$$
\operatorname{Pr}\left(\left|t_{k}\right|\left|q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right|>\frac{\nu}{9}\right) \leq\left(\frac{9}{\nu\left|t_{k}\right|}\right)^{r} \mathbb{E}\left|q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right|^{r} .
$$

Now, by Lemma 11, we have

$$
\mathbb{E}\left(\left.\left|q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right|^{r} \right\rvert\, D_{k}\right) \leq p^{-r} K_{r}\left(\left[\nu_{4} \operatorname{tr}\left(D_{k} D_{k}^{*}\right)\right]^{r / 2}+\nu_{2 r} \operatorname{tr}\left(D_{k} D_{k}^{*}\right)^{r / 2}\right)
$$

Since $\left\|D_{k}\right\| \leq 1 / v$, we have, for $z \in \mathcal{K}_{H}$,

$$
\mathbb{E}\left(\left.\left|q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right|^{r} \right\rvert\, D_{k}\right) \leq p^{-r} K_{r}\left(\nu_{4} \eta^{-r} p^{r / 2}+\nu_{2 r} p \eta^{-r}\right)
$$

Therefore,

$$
\operatorname{Pr}\left(\left|t_{k}\right|\left|q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right|>\frac{\nu}{9}\right) \leq\left(\frac{9}{\nu\left|t_{k}\right|}\right)^{r} p^{-r} K_{r}\left(\nu_{4} \eta^{-r} p^{r / 2}+\nu_{2 r} p \eta^{-r}\right),
$$

and $\left|t_{k}\right|\left|q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right| \leq \frac{\nu}{9}$ occur w.o.p.
For $z \in \mathcal{K}_{p V}$, we have

$$
\begin{aligned}
& \operatorname{Pr}\left(\left|t_{k}\right|\left|q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right|>\frac{\nu}{9}, Q_{p, k}\right) \\
= & \operatorname{Pr}\left(\left|t_{k}\right|\left|q_{k}^{\prime} D_{k} q_{k} \mathbf{1}_{Q_{p, k}}-\frac{1}{p} \operatorname{tr} D_{k} \mathbf{1}_{Q_{p, k}}\right|>\frac{\nu}{9}\right)
\end{aligned}
$$

and since the eigenvalues of $D_{k} \mathbf{1}_{Q_{p, k}}$ are bounded by $\eta^{-1}$, we get, using the above line of arguments,

$$
\operatorname{Pr}\left(\left|t_{k}\right|\left|q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right|>\frac{\nu}{9}, Q_{p, k}\right) \leq\left(\frac{9}{\nu\left|t_{k}\right|}\right)^{r} p^{-r} K_{r}\left(\nu_{4} \eta^{-r} p^{r / 2}+\nu_{2 r} p \eta^{-r}\right)
$$

Since $Q_{p, k}$ occur w.o.p., $\left|t_{k}\right|\left|q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right| \leq \frac{\nu}{9}$ occur w.o.p. too.

Bound on $\frac{1}{p}\left(\operatorname{tr} D_{k}-\operatorname{tr} D\right)$.
Inequalities (22) and (29) imply that

$$
\left|\frac{1}{p}\left(\operatorname{tr} D_{k}-\operatorname{tr} D\right)\right| \leq \frac{1}{p \eta}
$$

for $z \in \mathcal{K}_{H}$, and

$$
\left|\frac{1}{p}\left(\operatorname{tr} D_{k} \mathbf{1}_{Q_{p, a l l}}-\operatorname{tr} D \mathbf{1}_{Q_{p, a l l}}\right)\right|<\frac{1}{p} \frac{3}{2} \eta^{-1}
$$

for $z \in \mathcal{K}_{V}$. Therefore, events

$$
\left|\frac{\left|t_{k}\right|}{p}\left(\operatorname{tr} D_{k}-\operatorname{tr} D\right)\right| \leq \frac{\nu}{9}
$$

occur w.o.p.
To summarize, we have shown that $\left|t_{k}\right|\left|\varepsilon_{k}\right| \leq \nu / 3$ occur w.o.p. This implies that

$$
\left|\beta_{k}\right|>\nu / 3
$$

occur w.o.p. for any $z \in \mathcal{K}_{p}$.

## The deterministic term.

Let

$$
W_{p}=\cap_{k=1}^{n_{1}+n_{2}}\left\{\left|\beta_{k}\right|>\nu / 3\right\} \cap Q_{p, \text { all }}
$$

Note that $W_{p}$ occur w.o.p. This implies that $\sqrt{p}\left(\mathbb{E} \hat{m}\left(z ; l_{i}\right)-m_{p}\left(z ; l_{i}\right)\right) \rightarrow 0$ uniformly in $z \in \mathcal{K}_{p}$ if and only if $\sqrt{p}\left(\mathbb{E}\left[\hat{m}\left(z ; l_{i}\right) \mathbf{1}_{W_{p}}\right]-m_{p}\left(z ; l_{i}\right)\right) \rightarrow 0$ uniformly in $z \in \mathcal{K}_{p}$. Indeed, for $z \in \mathcal{K}_{p}$, $\left|\hat{m}\left(z ; l_{i}\right)\right| \leq 1 /|\operatorname{Im} z| \leq p^{-2}$, and

$$
\sqrt{p \mathbb{E}}\left|\hat{m}\left(z ; l_{i}\right) \mathbf{1}_{W_{p}^{c}}\right| \leq p^{-2} \sqrt{p} \operatorname{Pr}\left\{W_{p}^{c}\right\} \rightarrow 0
$$

Let us denote $\hat{m}\left(z ; l_{i}\right) \mathbf{1}_{W_{p}}$ as $\bar{m}\left(z ; l_{i}\right) \equiv \bar{m}(z)$ for brevity. Note that Lemma 10 implies that $\mathbb{E} \bar{m}(z)-$ $m(z) \rightarrow 0$ uniformly in $z \in \mathcal{K}_{p}$.

Consider

$$
\delta(z)=z \mathbb{E} \bar{m}(z)-\frac{c_{1}+c_{2}}{c_{1} c_{2}} \int \frac{t \mathbb{E} \bar{m}(z) \mathrm{d} \hat{\mathcal{T}}(t)}{1+t \mathbb{E} \bar{m}(z)}+1
$$

which is the left hand side of $(8)$ where $m_{p}(z)$ is replaced by $\mathbb{E} \bar{m}(z)$. We have

$$
\mathbb{E} \bar{m}(z)=-\left(z-\frac{c_{1}+c_{2}}{c_{1} c_{2}} \int \frac{t \mathrm{~d} \hat{\mathcal{T}}(t)}{1+t \mathbb{E} \bar{m}(z)}-\frac{\delta(z)}{\mathbb{E} \bar{m}(z)}\right)^{-1}
$$

Subtracting (7) and rearranging, we obtain

$$
\mathbb{E} \bar{m}(z)-m_{p}(z)=-\delta m_{p}\left[1-m_{p} \mathbb{E} \bar{m} \frac{c_{1}+c_{2}}{c_{1} c_{2}} \int \frac{t^{2} \mathrm{~d} \hat{\mathcal{T}}(t)}{(1+t \mathbb{E} \bar{m})\left(1+t m_{p}\right)}\right]^{-1}
$$

where we omit the dependence of $\delta, m_{p}$ and $\bar{m}$ on $z$ and $l_{i}$ to make the displayed formula easier to read. We will omit this dependence in what follows to make notations more compact.

To establish that $\sqrt{p}\left|\mathbb{E} \bar{m}-m_{p}\right| \rightarrow 0$, uniformly in $z \in \mathcal{K}_{p}$, it is sufficient to show that
(a) for sufficiently large $p,\left|m_{p}\right|$ is bounded, uniformly in $z \in \mathcal{K}_{p}$,
(b) for sufficiently large $p$, the absolute value of the square bracket in (30) is bounded away from zero, uniformly in $z \in \mathcal{K}_{p}$,
(c) $\sqrt{p}|\delta| \rightarrow 0$, uniformly in $z \in \mathcal{K}_{p}$.
(a) follows from the fact that $m_{p}$ converges to $m$, uniformly in $z \in \mathcal{K}_{p}$, which was established in Lemma 10.

To establish (b), it is sufficient to show that there exists $\xi \in(0,1)$ such that, for sufficiently large $p$,

$$
\begin{equation*}
|\mathbb{E} \bar{m}|^{2} \frac{c_{1}+c_{2}}{c_{1} c_{2}} \int \frac{t^{2} \mathrm{~d} \hat{\mathcal{T}}(t)}{|1+t \mathbb{E} \bar{m}|^{2}}<\xi \tag{30}
\end{equation*}
$$

and

$$
\begin{equation*}
\left|m_{p}\right|^{2} \frac{c_{1}+c_{2}}{c_{1} c_{2}} \int \frac{t^{2} \mathrm{~d} \hat{\mathcal{T}}(t)}{\left|1+t m_{p}\right|^{2}}<\xi \tag{31}
\end{equation*}
$$

uniformly in $z \in \mathcal{K}_{p}$. Indeed, then (b) follows by the Cauchy-Schwarz inequality.
In fact, since $m_{p}$ and $\mathbb{E} \bar{m}$ converge to $m$, uniformly in $z \in \mathcal{K}_{p}$, it is sufficient to prove that there exists $\xi \in(0,1)$

$$
\begin{equation*}
|m|^{2} \frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t^{2} \mathrm{~d} \mathcal{T}(t)}{|1+t m|^{2}}<\xi, \tag{32}
\end{equation*}
$$

uniformly in $z \in \mathcal{K}$. Using (5), we obtain

$$
\begin{align*}
|m|^{2} \frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t^{2} \mathrm{~d} \mathcal{T}(t)}{|1+t m|^{2}}= & \frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t^{2} \mathrm{~d} \mathcal{T}(t)}{|1+t m|^{2}} \\
& \times\left|z-\frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t \mathrm{~d} \mathcal{T}(t)}{1+t m}\right|^{-2} . \tag{33}
\end{align*}
$$

Using (5) again, we obtain for $z=u+i v$,

$$
\operatorname{Im} m=\left(v+\operatorname{Im} m \frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t^{2} \mathrm{~d} \mathcal{T}(t)}{|1+t m|^{2}}\right)\left|z-\frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t \mathrm{~d} \mathcal{T}(t)}{1+t m}\right|^{-2}
$$

Combining this with (33), we obtain

$$
\begin{aligned}
|m|^{2} \frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t^{2} \mathrm{~d} \mathcal{T}(t)}{|1+t m|^{2}}= & \operatorname{Im} m \frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t^{2} \mathrm{~d} \mathcal{T}(t)}{|1+t m|^{2}} \\
& \times\left(v+\operatorname{Im} m \frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t^{2} \mathrm{~d} \mathcal{T}(t)}{|1+t m|^{2}}\right)^{-1} .
\end{aligned}
$$

On the other hand, $\operatorname{Im} m=\int v \frac{\mathrm{~d} F^{\gamma, \mathcal{T}}}{|\lambda-z|^{2}}$. Therefore,

$$
\begin{aligned}
|m|^{2} \frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t^{2} \mathrm{~d} \mathcal{T}(t)}{|1+t m|^{2}}= & \int \frac{\mathrm{d} F^{\gamma, \mathcal{T}}}{|\lambda-z|^{2}} \frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t^{2} \mathrm{~d} \mathcal{T}(t)}{|1+t m|^{2}} \\
& \times\left(1+\int \frac{\mathrm{d} F^{\gamma, \mathcal{T}}}{|\lambda-z|^{2}} \frac{\gamma_{1}+\gamma_{2}}{\gamma_{1} \gamma_{2}} \int \frac{t^{2} \mathrm{~d} \mathcal{T}(t)}{|1+t m|^{2}}\right)^{-1}
\end{aligned}
$$

The right hand side of the above equality is smaller than one for any $z \in \mathcal{K}$. Since it is continuous on $z \in \mathcal{K}$, there exists $\xi \in(0,1)$ such that (32) holds, uniformly in $z \in \mathcal{K}$.

It remains to establish (c). It is sufficient to show that $\sqrt{p}|\bar{\delta}| \rightarrow 0$, where

$$
\bar{\delta}=\delta-\operatorname{Pr}\left\{W_{p}^{c}\right\}
$$

Define

$$
u=\frac{c_{1}+c_{2}}{c_{1} c_{2}} \int \frac{t \mathrm{~d} \hat{\mathcal{T}}(t)}{1+t \mathbb{E} \bar{m}}=\frac{1}{p} \sum_{k=1}^{n_{1}+n_{2}} \frac{t_{k}}{1+t_{k} \mathbb{E} \bar{m}}
$$

We have

$$
\begin{aligned}
\bar{\delta} & =\operatorname{Pr}\left\{W_{p}\right\}-(u-z) \mathbb{E} \bar{m}=\frac{1}{p} \operatorname{tr} \mathbb{E}\left[\left(A\left(l_{i}\right)-u I\right)\left(A\left(l_{i}\right)-z I\right)^{-1} \mathbf{1}_{W_{p}}\right] \\
& =\frac{1}{p} \operatorname{tr} \mathbb{E}\left[\sum_{k=1}^{n_{1}+n_{2}} t_{k} q_{k} q_{k}^{\prime} D \mathbf{1}_{W_{p}}-u D \mathbf{1}_{W_{p}}\right]
\end{aligned}
$$

Using the identity (16), we obtain

$$
\bar{\delta}=\frac{1}{p} \sum_{k=1}^{n_{1}+n_{2}} \mathbb{E}\left[\frac{t_{k} q_{k}^{\prime} D_{k} q_{k}}{\beta_{k}} \mathbf{1}_{W_{p}}\right]-y \mathbb{E} \bar{m}
$$

Using the definition of $u$, we can continue

$$
\bar{\delta}=\frac{1}{p} \sum_{k=1}^{n_{1}+n_{2}} \frac{t_{k}}{1+t_{k} \mathbb{E} \bar{m}} \mathbb{E}\left[\frac{t_{k} q_{k}^{\prime} D_{k} q_{k}}{\beta_{k}} \mathbf{1}_{W_{p}} \frac{1+t_{k} \mathbb{E} \bar{m}}{t_{k}}-\mathbb{E} \bar{m}\right]
$$

The term in the square brackets equals

$$
\begin{aligned}
& \frac{q_{k}^{\prime} D_{k} q_{k}}{\beta_{k}} \mathbf{1}_{W_{p}}+\left[\frac{t_{k} q_{k}^{\prime} D_{k} q_{k}}{\beta_{k}} \mathbf{1}_{W_{p}}-1\right] \mathbb{E} \bar{m} \\
= & \frac{q_{k}^{\prime} D_{k} q_{k}}{\beta_{k}} \mathbf{1}_{W_{p}}-\frac{1}{\beta_{k}} \mathbf{1}_{W_{p}} \mathbb{E} \bar{m}-\mathbf{1}_{W_{p}^{c}} \mathbb{E} \bar{m} \\
= & \frac{\varepsilon_{k}}{\beta_{k}} \mathbf{1}_{W_{p}}+\frac{\mathbb{E} \hat{m}-\mathbb{E} \bar{m}}{\beta_{k}} \mathbf{1}_{W_{p}}-\mathbf{1}_{W_{p}^{c}} \mathbb{E} \bar{m},
\end{aligned}
$$

where

$$
\varepsilon_{k}=q_{k}^{\prime} D_{k} q_{k}-\mathbb{E} \hat{m}
$$

Therefore, to establish (c), it is sufficient to show that

$$
\begin{gather*}
\left|\mathbb{E} \frac{\varepsilon_{k}}{\beta_{k}} \mathbf{1}_{W_{p}}\right|=O\left(p^{-1 / 2}\right),  \tag{34}\\
\left|\mathbb{E} \frac{\mathbb{E} \hat{m}-\mathbb{E} \bar{m}}{\beta_{k}} \mathbf{1}_{W_{p}}\right|=O\left(p^{-1 / 2}\right), \tag{35}
\end{gather*}
$$

and

$$
\begin{equation*}
\mathbb{E}\left|\mathbf{1}_{W_{p}^{c}} \mathbb{E} \bar{m}\right|=O\left(p^{-1 / 2}\right) \tag{36}
\end{equation*}
$$

uniformly in $z \in \mathcal{K}_{p}$. It is because $\max _{k}\left|\frac{t_{k}}{1+t_{k} \mathbb{E} \bar{m}}\right|$ is bounded uniformly in $z \in \mathcal{K}_{p}$, which follows from the uniform in $z \in \mathcal{K}_{p}$ convergence of $\mathbb{E} \hat{m}$ (and thus, of $\mathbb{E} \bar{m}$ ) to $m$, established by Lemma 10.

Equality (36) immediately follows from the fact that $W_{p}$ occur w.o.p. and from the boundedness of $\mathbb{E} \bar{m}$ (it converges to $m(z)$ which is bounded on $z \in \mathcal{K}$ ). For (35), we have

$$
\mathbb{E} \frac{\mathbb{E} \hat{m}-\mathbb{E} \bar{m}}{\beta_{k}} \mathbf{1}_{W_{p}}=\mathbb{E}\left(\hat{m} \mathbf{1}_{W_{p}^{c}}\right) \mathbb{E}\left(\beta_{k}^{-1} \mathbf{1}_{W_{p}}\right)
$$

But

$$
\begin{equation*}
\mathbb{E}\left(\beta_{k}^{-1} \mathbf{1}_{W_{p}}\right) \leq \frac{3}{\nu} . \tag{37}
\end{equation*}
$$

Further, $|\hat{m}| \leq 1 /|v| \leq p^{2}$ for $z \in \mathcal{K}_{p}$, and therefore,

$$
\begin{equation*}
\mathbb{E}\left(\hat{m} \mathbf{1}_{W_{p}^{c}}\right) \leq p^{2} \operatorname{Pr}\left\{W_{p}^{c}\right\} . \tag{38}
\end{equation*}
$$

Inequalities (37), (38), and the fact that $W_{p}$ occur w.o.p. imply (35).
The following lemma subsumes (34) by proving a stronger statement.
Lemma 12. There exists a constant $C$ such that, for sufficiently large $p$, for any $z \in \mathcal{K}_{p}$,

$$
\left|\mathbb{E} \frac{\varepsilon_{k}}{\beta_{k}} \mathbf{1}_{W_{p}}\right| \leq C p^{-1} .
$$

Proof: We need the following decomposition of $1 / \beta_{k}$ :

$$
\frac{1}{\beta_{k}}=\frac{1}{b_{k}}-\frac{1}{\beta_{k}} \frac{1}{b_{k}} \gamma_{k},
$$

where

$$
b_{k}=1+\frac{1}{p} t_{k} \mathbb{E} \operatorname{tr} D_{k}
$$

and

$$
\begin{aligned}
\gamma_{k} & =t_{k}\left[q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \mathbb{E} \operatorname{tr} D_{k}\right] \\
& =t_{k} \varepsilon_{k}+t_{k} \frac{1}{p} \mathbb{E}\left[\operatorname{tr} D-\operatorname{tr} D_{k}\right]
\end{aligned}
$$

Using the decomposition, we obtain

$$
\left|\mathbb{E} \frac{\varepsilon_{k}}{\beta_{k}} \mathbf{1}_{W_{p}}\right| \leq\left|\mathbb{E} \frac{\varepsilon_{k}}{b_{k}} \mathbf{1}_{W_{p}}\right|+\left|\mathbb{E} \frac{t_{k} \varepsilon_{k}^{2}}{\beta_{k} b_{k}} \mathbf{1}_{W_{p}}\right|+\left|\frac{1}{p} \mathbb{E}\left[\operatorname{tr} D-\operatorname{tr} D_{k}\right]\right|\left|\mathbb{E} \frac{t_{k} \varepsilon_{k}}{\beta_{k} b_{k}} \mathbf{1}_{W_{p}}\right| .
$$

Bound on $\left|\mathbb{E} \frac{\varepsilon_{k}}{b_{k}} \mathbf{1}_{W_{p}}\right|$.
We have

$$
\mathbb{E}\left[\frac{\varepsilon_{k}}{b_{k}} \mathbf{1}_{W_{p}}\right]=\frac{1}{b_{k}} \mathbb{E} \varepsilon_{k}-\frac{1}{b_{k}} \mathbb{E}\left[\varepsilon_{k} \mathbf{1}_{W_{p}^{c}}\right] .
$$

Since

$$
\mathbb{E}\left(\varepsilon_{k} \mid D_{k}\right)=\frac{1}{p}\left(\operatorname{tr} D_{k}-\mathbb{E} \operatorname{tr} D\right)
$$

we have

$$
\begin{equation*}
\mathbb{E}\left(\varepsilon_{k}\right)=\frac{1}{p} \mathbb{E}\left(\operatorname{tr} D_{k}-\operatorname{tr} D\right)=-\frac{1}{p} \mathbb{E}\left(y-y_{k}\right) \tag{39}
\end{equation*}
$$

By (22), we have

$$
\begin{equation*}
\frac{1}{p}\left|\mathbb{E}\left(y-y_{k}\right)\right|<\frac{1}{p \eta} \text { for } z \in \mathcal{K}_{H}, \tag{40}
\end{equation*}
$$

and by (29), we have

$$
\begin{align*}
\frac{1}{p}\left|\mathbb{E}\left(y-y_{k}\right)\right| & <\frac{1}{p} \frac{3}{2} \eta+\frac{1}{p}\left|\mathbb{E}\left(\operatorname{tr} D_{k} \mathbf{1}_{W_{p}^{c}}-\operatorname{tr} D \mathbf{1}_{W_{p}^{c}}\right)\right|  \tag{41}\\
& \leq \frac{1}{p} \frac{3}{2} \eta+\frac{1}{p}\left(2 p^{3}\right) \operatorname{Pr}\left(W_{p}^{c}\right) \leq \frac{1}{p} 2 \eta^{-1}
\end{align*}
$$

for sufficiently large $p$ and $z \in \mathcal{K}_{p V}$. Using (40) and (41) in (39), we obtain

$$
\begin{equation*}
\left|\mathbb{E}\left(\varepsilon_{k}\right)\right| \leq \frac{1}{p} \frac{3}{\eta} \tag{42}
\end{equation*}
$$

for any $z \in \mathcal{K}_{p}$.
Further, since the eigenvalues of $D_{k}$ are no larger than $1 /|v|$ by absolute value, and since $1 /|v| \leq$ $p^{2}$ for $z \in \mathcal{K}_{p}$, we have

$$
\begin{aligned}
\left|\mathbb{E}\left[\varepsilon_{k} \mathbf{1}_{W_{p}^{c}}\right]\right| & \leq p^{2}\left|\mathbb{E}\left[q_{k}^{\prime} q_{k} \mathbf{1}_{W_{p}^{c}}\right]\right|+\mathbb{E} \hat{m} \operatorname{Pr}\left(W_{p}^{c}\right) \\
& \leq p^{2}\left(\mathbb{E}\left[\left(q_{k}^{\prime} q_{k}\right)^{2}\right]\right)^{1 / 2} \operatorname{Pr}\left(W_{p}^{c}\right)^{1 / 2}+\mathbb{E} \hat{m} \operatorname{Pr}\left(W_{p}^{c}\right) .
\end{aligned}
$$

Therefore,

$$
\begin{equation*}
\left|\mathbb{E}\left[\varepsilon_{k} \mathbf{1}_{W_{p}^{c}}\right]\right| \leq p^{-1} C \tag{43}
\end{equation*}
$$

for sufficiently large $p$ and some $C$.
Next, by Lemma $10, b_{k}=1+\frac{1}{p} t_{k} \mathbb{E} \operatorname{tr} D_{k} \rightarrow 1+t_{k} m(z)$ and hence,

$$
\begin{equation*}
\left|b_{k}\right|>\nu / 2 \tag{44}
\end{equation*}
$$

for sufficiently large $p$. Combining (42), (43), and (44), we obtain

$$
\begin{equation*}
\left|\mathbb{E} \frac{\varepsilon_{k}}{b_{k}} \boldsymbol{1}_{W_{p}}\right| \leq p^{-1} C \tag{45}
\end{equation*}
$$

for sufficiently large $p$ and some $C$.
Bound on $\left|\mathbb{E} \frac{t_{k} \varepsilon_{k}^{2}}{\beta_{k} b_{k}} \mathbf{1}_{W_{p}}\right|$.
By (44), for sufficiently large $p$, we have

$$
\begin{equation*}
\left|\mathbb{E} \frac{t_{k} \varepsilon_{k}^{2}}{\beta_{k} b_{k}} \mathbf{1}_{W_{p}}\right| \leq \frac{6 t_{k}}{\nu^{2}} \mathbb{E}\left(\varepsilon_{k}^{2} \mathbf{1}_{W_{p}}\right) . \tag{46}
\end{equation*}
$$

Consider the decomposition

$$
\begin{equation*}
\varepsilon_{k}=\left(q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right)+\frac{1}{p}\left(\operatorname{tr} D_{k}-\mathbb{E} \operatorname{tr} D\right) \tag{47}
\end{equation*}
$$

By Lemma 11,

$$
\mathbb{E}\left(\left.\left(q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right)^{2} \right\rvert\, D_{k}\right) \leq \frac{1}{p^{2}} K_{2}\left(\left[\nu_{4} \operatorname{tr}\left(D_{k} D_{k}^{*}\right)\right]^{1 / 2}+\nu_{4} \operatorname{tr}\left(D_{k} D_{k}^{*}\right)\right)
$$

Since the eigenvalues of $D_{k}$ are bounded by absolute value by

$$
c_{z} \equiv \min \left\{\frac{1}{|v|}, \eta^{-1} \mathbf{1}_{Q_{p, a l l}} \mathbf{1}_{z \in \mathcal{K}_{V}}\right\},
$$

we have

$$
\mathbb{E}\left(\left.\left(q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right)^{2} \right\rvert\, D_{k}\right) \leq \frac{1}{p^{2}} K_{2}\left(\left[\nu_{4} p c_{z}^{2}\right]^{1 / 2}+\nu_{4} p c_{z}^{2}\right)
$$

and

$$
\begin{aligned}
\mathbb{E}\left(q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right)^{2} \leq & \frac{1}{p^{2}} K_{2}\left(\nu_{4} p\right)^{1 / 2} p^{2} \operatorname{Pr}\left\{Q_{p, \text { all }}^{c}\right\}+\frac{1}{p^{2}} K_{2}\left(\nu_{4} p\right)^{1 / 2} \eta^{-1} \\
& +\frac{1}{p^{2}} K_{2} \nu_{4} p p^{4} \operatorname{Pr}\left\{Q_{p, \text { all }}^{c}\right\}+\frac{1}{p^{2}} K_{2} \nu_{4} p \eta^{-2}
\end{aligned}
$$

Therefore, for sufficiently large $p$,

$$
\begin{equation*}
\mathbb{E}\left(q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right)^{2} \leq C p^{-1} \tag{48}
\end{equation*}
$$

for some $C$.
For the second part of the decomposition (47), we have

$$
\begin{equation*}
\mathbb{E} \frac{1}{p^{2}}\left(\operatorname{tr} D_{k}-\mathbb{E} \operatorname{tr} D\right)^{2} \leq \mathbb{E} \frac{2}{p^{2}}\left[\left(\operatorname{tr} D_{k}-\mathbb{E} \operatorname{tr} D_{k}\right)^{2}+\left(\mathbb{E} \operatorname{tr} D_{k}-\mathbb{E} \operatorname{tr} D\right)^{2}\right] \tag{49}
\end{equation*}
$$

Note that $\left(\mathbb{E}\left[\operatorname{tr} D_{k}-\operatorname{tr} D\right]\right)^{2}$ is bounded by (40) and (41). Let us now prove that

$$
\mathbb{E}\left(\operatorname{tr} D_{k}-\mathbb{E} \operatorname{tr} D_{k}\right)^{2} \leq C p
$$

for some $C$. We will prove this inequality for $D_{k}$ replaced by $D$ to ease notation. The proof for $D_{k}$ is very similar.

Recall that

$$
\operatorname{tr} D-\mathbb{E} \operatorname{tr} D=\sum_{k=1}^{n_{1}+n_{2}}\left(\mathbb{E}_{k}-\mathbb{E}_{k-1}\right)\left(y-y_{k}\right)
$$

and

$$
\mathbb{E}(\operatorname{tr} D-\mathbb{E} \operatorname{tr} D)^{2}=\sum_{k=1}^{n_{1}+n_{2}} \mathbb{E} z_{k}^{2}
$$

where $z_{k}=\left(\mathbb{E}_{k}-\mathbb{E}_{k-1}\right)\left(y-y_{k}\right)$.

Using (22) and (29), we obtain

$$
\begin{aligned}
\mathbb{E} z_{k}^{2}= & \mathbb{E}\left[\mathbb{E}_{k}\left(y-y_{k}\right)-\mathbb{E}_{k-1}\left(y-y_{k}\right)\right]^{2} \leq \\
& 2 \mathbb{E}\left[\mathbb{E}_{k}\left(y-y_{k}\right)\right]^{2}+2 \mathbb{E}\left[\mathbb{E}_{k-1}\left(y-y_{k}\right)\right]^{2} \\
\leq & 2 \mathbb{E} \mathbb{E}_{k}\left(y-y_{k}\right)^{2}+2 \mathbb{E}_{\mathbb{E}_{k-1}}\left(y-y_{k}\right)^{2} \\
= & 4 \mathbb{E}\left(y-y_{k}\right)^{2} \\
\leq & 4 p^{4} \operatorname{Pr}\left\{Q_{p, \text { all }}^{c}\right\}+4\left\{\frac{1}{\eta^{2}}+\left(\frac{3}{2 \eta}\right)^{2}\right\}
\end{aligned}
$$

Hence, $\mathbb{E}(\operatorname{tr} D-\mathbb{E} \operatorname{tr} D)^{2} \leq C p$ (and $\left.\mathbb{E}\left(\operatorname{tr} D_{k}-\mathbb{E} \operatorname{tr} D_{k}\right)^{2} \leq C p\right)$ for some $C$, and, from (49), we have

$$
\begin{equation*}
\mathbb{E} \frac{1}{p^{2}}\left(\operatorname{tr} D_{k}-\mathbb{E} \operatorname{tr} D\right)^{2} \leq C p^{-1} \tag{50}
\end{equation*}
$$

for some $C$. Now, (47), (48), and (50) imply that

$$
\begin{equation*}
\mathbb{E}\left(\varepsilon_{k}^{2}\right) \leq C p^{-1} \tag{51}
\end{equation*}
$$

for some $C$. Therefore, by (46),

$$
\begin{equation*}
\left|\mathbb{E} \frac{t_{k} \varepsilon_{k}^{2}}{\beta_{k} b_{k}} \mathbf{1}_{W_{p}}\right| \leq C p^{-1} \tag{52}
\end{equation*}
$$

for some $C$.
Bound on $\left|\frac{1}{p} \mathbb{E}\left[\operatorname{tr} D-\operatorname{tr} D_{k}\right]\right|\left|\mathbb{E} \frac{t_{k} \varepsilon_{k}}{\beta_{k} b_{k}} \mathbf{1}_{W_{p}}\right|$.
By (40) and (41),

$$
\begin{equation*}
\left|\frac{1}{p} \mathbb{E}\left[\operatorname{tr} D-\operatorname{tr} D_{k}\right]\right| \leq p^{-1} \frac{3}{\eta} \tag{53}
\end{equation*}
$$

Further, by the Cauchy-Schwarz inequality,

$$
\left|\mathbb{E} \frac{t_{k} \varepsilon_{k}}{\beta_{k} b_{k}} \mathbf{1}_{W_{p}}\right|^{2} \leq \mathbb{E} \varepsilon_{k}^{2} \mathbb{E}\left(\frac{t_{k}^{2}}{\beta_{k}^{2} b_{k}^{2}} \mathbf{1}_{W_{p}}\right)
$$

Inequality (51) and the boundedness of $\beta_{k}^{2} b_{k}^{2}$ away from zero on $W_{p}$ imply that the right hand side of the above inequality is bounded. From this and (53) we see that

$$
\begin{equation*}
\left|\frac{1}{p} \mathbb{E}\left[\operatorname{tr} D-\operatorname{tr} D_{k}\right]\right|\left|\mathbb{E} \frac{t_{k} \varepsilon_{k}}{\beta_{k} b_{k}} \mathbf{1}_{W_{p}}\right| \leq C p^{-1} \tag{54}
\end{equation*}
$$

for some $C$.
The Lemma follows from (45), (52), and (54).
The validity of (36) is implied by the validity of the Lemma, and thus, $\sqrt{p}|\bar{\delta}| \rightarrow 0$. This concludes our proof of the deterministic term's convergence (15).

## The stochastic term.

To get the correct order of magnitude for the fluctuations of

$$
z_{k}=\left(\mathbb{E}_{k}-\mathbb{E}_{k-1}\right)\left(y-y_{k}\right),
$$

we need a finer decomposition. First, define some conditional means and residuals:

$$
\begin{gathered}
\beta_{k}^{0}=\mathbb{E}\left(\beta_{k} \mid D_{k}\right)=1+\frac{1}{p} t_{k} \operatorname{tr} D_{k}, \\
\gamma_{k}^{0}=\mathbb{E}\left(\gamma_{k} \mid D_{k}\right)=-\frac{1}{p} t_{k} \operatorname{tr} D_{k}^{2}, \\
\varepsilon_{1 k}=\beta_{k}-\beta_{k}^{0}=t_{k}\left(q_{k}^{\prime} D_{k} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}\right), \\
\varepsilon_{2 k}=\gamma_{k}-\gamma_{k}^{0}=-t_{k}\left(q_{k}^{\prime} D_{k}^{2} q_{k}-\frac{1}{p} \operatorname{tr} D_{k}^{2}\right) .
\end{gathered}
$$

Rewrite $y-y_{k}$ in terms of these means and residulas

$$
\begin{aligned}
y-y_{k} & =\frac{\gamma_{k}}{\beta_{k}}=\frac{\gamma_{k}^{0}}{\beta_{k}^{0}}-\frac{\varepsilon_{1 k} \gamma_{k}}{\beta_{k}^{0} \beta_{k}}+\frac{\varepsilon_{2 k}}{\beta_{k}^{0}} \\
& =r_{0 k}+r_{1 k}+r_{2 k}
\end{aligned}
$$

The integrable terms ${ }^{1} r_{0 k}$ do not contain any variables from $\mathcal{F}_{k} \backslash \mathcal{F}_{k-1}$. Therefore,

$$
\left(\mathbb{E}_{k}-\mathbb{E}_{k-1}\right) r_{0 k}=0 .
$$

Thus, they disappear from the martingale differences and

$$
z_{k}=\left(\mathbb{E}_{k}-\mathbb{E}_{k-1}\right)\left(y-y_{k}\right)=\left(\mathbb{E}_{k}-\mathbb{E}_{k-1}\right)\left(r_{1 k}+r_{2 k}\right)
$$

## Bounds on $r_{m k}$.

Crude bounds. By (22),

$$
\left|\gamma_{k} / \beta_{k}\right| \leq 1 /|v| .
$$

Further

$$
\begin{aligned}
\left|\frac{\gamma_{k}^{0}}{\beta_{k}^{0}}\right| & \leq \frac{\frac{1}{p}\left|t_{k}\right| \sum_{j=1}^{p}\left|\mu_{j k}-z\right|^{-2}}{\left|\operatorname{Im} \beta_{k}^{0}\right|} \\
& =\frac{\frac{1}{p}\left|t_{k}\right| \sum_{j=1}^{p}\left|\mu_{j k}-z\right|^{-2}}{\frac{1}{p}\left|t_{k}\right||v| \sum_{j=1}^{p}\left|\mu_{j k}-z\right|^{-2}}=\frac{1}{|v|}
\end{aligned}
$$

Next,

$$
\left|\beta_{k}^{0}\right|>\frac{1}{p}\left|t_{k}\right|\left|\operatorname{Im} D_{k}\right|=\frac{1}{p}\left|t_{k}\right| \sum_{j=1}^{p} \frac{v}{\left|\mu_{j k}-z\right|^{2}},
$$

whereas

$$
\left|\varepsilon_{2 k}\right| \leq\left|t_{k}\right| \sum_{j=1}^{p} \frac{\left(q_{k}^{\prime} e_{j}\right)^{2}+1 / p}{\left|\mu_{j k}-z\right|^{2}} \leq\left|t_{k}\right|\left(q_{k}^{\prime} q_{k}+1 / p\right) \sum_{j=1}^{p} \frac{1}{\left|\mu_{j k}-z\right|^{2}}
$$

$$
\begin{aligned}
& { }^{1}\left|\gamma_{k}^{0} / \beta_{k}^{0}\right| \leq \frac{1}{p} t_{k} \sum_{j=1}^{p}\left|\mu_{j k}-z\right|^{-2} /\left(\left|\operatorname{Im} \beta_{k}^{0}\right|\right) \\
& \qquad \operatorname{Im} \beta_{k}^{0}=\frac{1}{p} t_{k} v \sum_{j=1}^{p}\left|\mu_{j k}-z\right|^{-2}
\end{aligned}
$$

therefore, $\left|\gamma_{k}^{0} / \beta_{k}^{0}\right| \leq 1 / v$, and thus, $\gamma_{k}^{0} / \beta_{k}^{0}$ is integrable.

Therefore,

$$
\left|r_{2 k}\right| \leq \frac{\left|t_{k}\right|\left(q_{k}^{\prime} q_{k}+1 / p\right) \sum_{j=1}^{p} \frac{1}{\left|\mu_{j k}-z\right|^{2}}}{\frac{1}{p}\left|t_{k}\right| \sum_{j=1}^{p} \frac{v}{\left|\mu_{j k}-z\right|^{2}}}=\frac{1}{v}\left(p q_{k}^{\prime} q_{k}+1\right)
$$

Finally,

$$
\left|r_{1 k}\right| \leq\left|\frac{\gamma_{k}}{\beta_{k}}\right|+\left|r_{0 k}\right|+\left|r_{2 k}\right| \leq \frac{1}{v}\left(p q_{k}^{\prime} q_{k}+3\right)
$$

Fine bounds.
Consider the event

$$
B_{k}=\left\{\left|\beta_{k}\right|>\nu / 3,\left|\beta_{k}^{0}\right|>\nu / 3, p q_{k}^{\prime} q_{k} \leq 2 p\right\} \cap Q_{p, \text { all }}
$$

On this event, which occurs w.o.p. (a proof of the fact that $\left|\beta_{k}^{0}\right|>\nu / 3$ occur w.o.p. is actually contained in the above proof of the fact that $\left|\beta_{k}\right|>\nu / 3$ occur w.o.p.), we have

$$
\left|\gamma_{k}\right| \leq\left|t_{k}\right| 2 \eta
$$

for $z \in \mathcal{K}_{p}$ (note that the distance from any $\mu_{j k}$ to $\mathcal{K}$ is bounded by $\eta$ on $Q_{p, \text { all }}$ ). Therefore, on $B_{k}$, we have

$$
\left|r_{1 k}\right| \leq\left|\varepsilon_{1 k}\right| \frac{9\left|t_{k}\right| 2}{\nu^{2} \eta}
$$

and

$$
\left|r_{2 k}\right| \leq\left|\varepsilon_{1 k}\right| 3 / \nu
$$

Bounds on $\varepsilon_{1 k}$ and $\varepsilon_{2 k}$.
We have

$$
\mathbb{E}_{k-1}\left[\left|\varepsilon_{1 k}\right|^{2 l} \mathbf{1}_{Q_{p k}}\right]=p^{-2 l} \mathbb{E}_{k-1}\left[\mathbf{1}_{Q_{p k}} \mathbb{E}\left(\left|p q_{k}^{\prime} D_{k} q_{k}-\operatorname{tr} D_{k}\right|^{2 l} \mid D_{k}\right)\right]
$$

Using Lemma 11, we continue, for $l \geq 1$,

$$
\begin{aligned}
\mathbb{E}_{k-1}\left[\left|\varepsilon_{1 k}\right|^{2 l} \mathbf{1}_{Q_{p k}}\right] & \leq p^{-2 l} \mathbb{E}_{k-1}\left[\mathbf{1}_{Q_{p k}} K_{2 l}\left(\nu_{4}^{l} p^{l} \eta^{-2 l}+\nu_{4 l} p \eta^{-2 l}\right)\right] \\
& \leq \tilde{K}_{l} p^{-l} \eta^{-2 l}
\end{aligned}
$$

Similarly,

$$
\mathbb{E}_{k-1}\left[\left|\varepsilon_{2 k}\right|^{2 l} \mathbf{1}_{Q_{p k}}\right] \leq \tilde{K}_{l} p^{-l} \eta^{-4 l}
$$

Let us prove the following lemma.
Lemma 13. For any $l \geq 1, \mathbb{E}\left|z_{k}\right|^{2 l} \leq C p^{-l}$ uniformly in $z \in \mathcal{K}_{p}$.
Proof: Set $W=r_{1 k}+r_{2 k}$ and observe that

$$
\begin{aligned}
\mathbb{E}_{k-1}\left|z_{k}\right|^{2 l} & =\mathbb{E}_{k-1}\left|\mathbb{E}_{k} W-\mathbb{E}_{k-1} W\right|^{2 l} \leq c_{2 l}\left(\mathbb{E}_{k-1}\left|\mathbb{E}_{k} W\right|^{2 l}+\mathbb{E}_{k-1}\left|\mathbb{E}_{k-1} W\right|^{2 l}\right) \\
& \leq c_{2 l}\left(\mathbb{E}_{k-1} \mathbb{E}_{k}|W|^{2 l}+\mathbb{E}_{k-1} \mathbb{E}_{k-1}|W|^{2 l}\right)=2 c_{2 l} \mathbb{E}_{k-1}|W|^{2 l} \\
& \leq c_{l}^{\prime}\left(\mathbb{E}_{k-1}\left|r_{1 k}\right|^{2 l}+\mathbb{E}_{k-1}\left|r_{2 k}\right|^{2 l}\right)
\end{aligned}
$$

using $(a+b)^{p} \leq c_{p}\left(a^{p}+b^{p}\right)$ for $c_{p}=2^{p-1}$ and $p \geq 2$.
Further

$$
\begin{aligned}
\mathbb{E}_{k-1}\left|r_{1 k}\right|^{2 l} & =\mathbb{E}_{k-1}\left(\left|r_{1 k}\right|^{2 l} \mathbf{1}_{B_{k}}\right)+\mathbb{E}_{k-1}\left(\left|r_{1 k}\right|^{2 l} \mathbf{1}_{B_{k}^{c}}\right) \\
& \leq\left(\frac{9\left|t_{k}\right| 2}{\nu^{2} \eta}\right)^{2 l} \mathbb{E}_{k-1}\left(\left|\varepsilon_{1 k}\right|^{2 l} \mathbf{1}_{B_{k}}\right)+\mathbb{E}_{k-1}\left(\left|r_{1 k}\right|^{2 l} \mathbf{1}_{B_{k}^{c}}\right) \\
& \leq\left(\frac{9\left|t_{k}\right| 2}{\nu^{2} \eta}\right)^{2 l} \mathbb{E}_{k-1}\left(\left|\varepsilon_{1 k}\right|^{2 l} \mathbf{1}_{Q_{p k}}\right)+\mathbb{E}_{k-1}\left(\left|r_{1 k}\right|^{2 l} \mathbf{1}_{B_{k}^{c}}\right)
\end{aligned}
$$

Taking unconditional expectations, we get

$$
\begin{aligned}
\mathbb{E}\left|r_{1 k}\right|^{2 l} & \leq\left(\frac{9\left|t_{k}\right| 2}{\nu^{2} \eta}\right)^{2 l} \mathbb{E}\left(\left|\varepsilon_{1 k}\right|^{2 l} \mathbf{1}_{Q_{p k}}\right)+\mathbb{E}\left(\left|r_{1 k}\right|^{2 l} \mathbf{1}_{B_{k}^{c}}\right) \\
& \leq\left(\frac{9\left|t_{k}\right| 2}{\nu^{2} \eta}\right)^{2 l} \tilde{K}_{l} p^{-l} \eta^{-2 l}+\mathbb{E}\left(\frac{1}{v^{2 l}}\left|p q_{k}^{\prime} q_{k}+3\right|^{2 l} \mathbf{1}_{B_{k}^{c}}\right) \\
& \leq\left(\frac{9\left|t_{k}\right| 2}{\nu^{2} \eta}\right)^{2 l} \tilde{K}_{l} p^{-l} \eta^{-2 l}+\frac{1}{v^{2 l}} \mathbb{E}\left(\left|p q_{k}^{\prime} q_{k}+3\right|^{4 l}\right)^{1 / 2}\left[\operatorname{Pr}\left(B_{k}^{c}\right)\right]^{1 / 2} \\
& \leq C p^{-l}
\end{aligned}
$$

A similar argument shows that

$$
\mathbb{E}\left|r_{2 k}\right|^{2 l} \leq C p^{-l} .
$$

To establish the stochastic bounds (14), we need to show that $\mathbb{E}|y-\mathbb{E} y|^{2 l} \leq C_{l}$ uniformly in $z \in \mathcal{K}_{p}$. By Rosenthal's inequality (see Theorem 2.12 of Hall and Heyde), we have for any $l \geq 1$,

$$
\mathbb{E}|y-\mathbb{E} y|^{2 l} \leq C_{l} \mathbb{E}\left[\left(\sum_{k=1}^{n_{1}+n_{2}} \mathbb{E}_{k-1} z_{k}^{2}\right)^{l}\right]+C_{l} \sum_{k=1}^{n_{1}+n_{2}} \mathbb{E}\left|z_{k}\right|^{2 l}
$$

That the second sum on the right is uniformly bounded follows immediately from Lemma 13.
Turn to the first sum. First, obtain the bound

$$
\left(\sum_{k=1}^{n_{1}+n_{2}} \mathbb{E}_{k-1} z_{k}^{2}\right)^{l} \leq c_{l}\left(R_{1 p}^{l}+R_{2 p}^{l}\right)
$$

where

$$
R_{m p}=\sum_{k=1}^{n_{1}+n_{2}} \mathbb{E}_{k-1}\left|r_{m k}\right|^{2}
$$

Recall that

$$
\begin{aligned}
\mathbb{E}_{k-1}\left|r_{1 k}\right|^{2} & \leq\left(\frac{9\left|t_{k}\right| 2}{\nu^{2} \eta}\right)^{2} \mathbb{E}_{k-1}\left(\left|\varepsilon_{1 k}\right|^{2} \mathbf{1}_{Q_{p k}}\right)+\mathbb{E}_{k-1}\left(\left|r_{1 k}\right|^{2} \mathbf{1}_{B_{k}^{c}}\right) \\
& \leq\left(\frac{9\left|t_{k}\right| 2}{\nu^{2} \eta}\right)^{2} \tilde{K}_{1} p^{-1} \eta^{-2}+\mathbb{E}_{k-1}\left(\left|r_{1 k}\right|^{2} \mathbf{1}_{B_{k}^{c}}\right)
\end{aligned}
$$

Therefore,

$$
R_{1 p} \leq C+\sum_{k=1}^{n_{1}+n_{2}} \mathbb{E}_{k-1}\left(\left|r_{1 k}\right|^{2} \mathbf{1}_{B_{k}^{c}}\right)
$$

and

$$
R_{1 p}^{l} \leq c_{l}\left[C^{l}+\left\{\sum_{k=1}^{n_{1}+n_{2}} \mathbb{E}_{k-1}\left(\left|r_{1 k}\right|^{2} \mathbf{1}_{B_{k}^{c}}\right)\right\}^{l}\right]
$$

On the other hand, for $l \geq 2$

$$
\begin{aligned}
\left\{\sum_{k=1}^{n_{1}+n_{2}} \mathbb{E}_{k-1}\left(\left|r_{1 k}\right|^{2} \mathbf{1}_{B_{k}^{c}}\right)\right\}^{l} & \leq\left(n_{1}+n_{2}\right)^{l-1} \sum_{k=1}^{n_{1}+n_{2}}\left[\mathbb{E}_{k-1}\left(\left|r_{1 k}\right|^{2} \mathbf{1}_{B_{k}^{c}}\right)\right]^{l} \\
& \leq\left(n_{1}+n_{2}\right)^{l-1} \sum_{k=1}^{n_{1}+n_{2}} \mathbb{E}_{k-1}\left(\left|r_{1 k}\right|^{2 l}\right) \mathbb{E}_{k-1} \mathbf{1}_{B_{k}^{c}} \\
& \leq\left(n_{1}+n_{2}\right)^{l-1} \sum_{k=1}^{n_{1}+n_{2}} \mathbb{E}_{k-1}\left(\left|\frac{1}{v}\left(p q_{k}^{\prime} q_{k}+3\right)\right|^{2 l}\right) \mathbb{E}_{k-1} \mathbf{1}_{B_{k}} \\
& \leq p^{K} \sum_{k=1}^{n_{1}+n_{2}} \mathbb{E}_{k-1} \mathbf{1}_{B_{k}^{c}}
\end{aligned}
$$

for some $K$. Taking unconditional expectations, we obtain the boundedness of $R_{1 p}^{l}$. The boundedness of $R_{2 p}^{l}$ is established similarly. This completes the proof of (14) and hence, of Lemma DJO6.
4.3. Proof of Lemma DJO7 (a CLT for quadratic forms). We will need the following two lemmas.

Lemma 14. (McLeish 1974) Let $\left\{X_{p r}, \mathcal{G}_{p r}, r=1, \ldots, p\right\}$ be a martingale difference array on the probability triple $(\Omega, \mathcal{G}, P)$. If the following conditions are satisfied: a) Lindeberg's condition: for all $\varepsilon>0, \sum_{r} \int_{\left|X_{p r}\right|>\varepsilon} X_{p r}^{2} \mathrm{~d} P \rightarrow 0$ as $\left.p \rightarrow \infty ; b\right) \sum_{r} X_{p r}^{2} \xrightarrow{\mathrm{P}} 1$, then $\sum_{r} X_{p r} \xrightarrow{d} N(0,1)$.

Proof: This is a consequence of Theorem (2.3) of McLeish (1974). Two conditions of the theorem: i) $\max _{r \leq p}\left|X_{p r}\right|$ is uniformly bounded in $L_{2}$ norm, and ii) $\max _{r \leq p}\left|X_{p r}\right| \xrightarrow{\mathrm{P}} 0$, are replaced here by the Lindeberg condition.

Lemma 15. (Hall and Heyde) Let $\left\{X_{p r}, \mathcal{G}_{p r}, r=1, \ldots, p\right\}$ be a martingale difference array, and define $V_{p j}^{2}=\sum_{r=1}^{j} E\left(X_{p r}^{2} \mid \mathcal{G}_{p, r-1}\right)$ and $U_{p j}^{2}=\sum_{r=1}^{j} X_{p r}^{2}$ for $r=1, \ldots, p$. Suppose that the conditional variances $V_{p p}^{2}$ are tight, that is $\sup _{p} P\left(V_{p p}^{2}>\varepsilon\right) \rightarrow 0$ as $\varepsilon \rightarrow \infty$, and that the conditional Lindeberg condition holds, that is, for all $\varepsilon>0, \sum_{r} E\left[X_{p r}^{2} \mathbf{1}\left\{\left|X_{p r}\right|>\varepsilon\right\} \mid \mathcal{G}_{p, r-1}\right] \xrightarrow{\mathrm{P}} 0$. Then $\max _{j}\left|U_{p j}^{2}-V_{p j}^{2}\right| \xrightarrow{\mathrm{P}} 0$.

Proof: This is a shortened version of Theorem 2.23 in Hall and Heyde (1980).
Let $f_{j}(\lambda), j=1, \ldots, J$, be such that $f_{j}(\lambda)=g_{j}(\lambda)$ for $\lambda \in\left[b_{i}, B_{i}\right]$ and $f_{j}(\lambda)=0$ otherwise. Consider random variables

$$
X_{p r}=\frac{1}{\sqrt{p}} \sum_{(j, s, t) \in \Theta} \gamma_{j s t} f_{j}\left(\mu_{r, i}\right)\left(\zeta_{r s} \zeta_{r t}-\delta_{s t}\right)
$$

where $\gamma_{j s t}$ are some constants. Let $\mathcal{G}_{p R}$ be the $\sigma$-algebra generated by $\mu_{1, i}, \ldots, \mu_{p, i}$ and $\zeta_{r s}$ with $r=1, \ldots, R ; s=1, \ldots, q$. Clearly, $\left\{X_{p r}, \mathcal{G}_{p r}, r=1, \ldots, p\right\}$ form a martingale difference array. Let
$K$ be the number of different triples $(j, s, t) \in \Theta$. Consider an arbitrary order in $\Theta$. In Hölder's inequality

$$
\sum_{a=1}^{K} y_{a} z_{a} \leq\left(\sum_{a=1}^{K}\left(y_{a}\right)^{b}\right)^{1 / b}\left(\sum_{a=1}^{K}\left(z_{a}\right)^{c}\right)^{1 / c}
$$

which holds for $y_{a}>0, z_{a}>0, b>1, c>1$, and $1 / b+1 / c=1$, take

$$
y_{a}=\left|\frac{1}{\sqrt{p}} \gamma_{j s t} f_{j}\left(\mu_{r, i}\right)\left(\zeta_{r s} \zeta_{r t}-\delta_{s t}\right)\right|,
$$

where $(j, s, t)$ is the $a$-th triple in $\Theta, z_{a}=1$, and $b=2+\delta$ for some $\delta>0$. Then, the inequality implies that

$$
\begin{equation*}
\left|X_{p r}\right|^{2+\delta} \leq K^{1+\delta} \rho_{i}^{2+\delta} \sum_{(j, s, t) \in \Theta}\left|\frac{1}{\sqrt{p}} \gamma_{j s t}\left(\zeta_{r s} \zeta_{r t}-\delta_{s t}\right)\right|^{2+\delta} \tag{55}
\end{equation*}
$$

where

$$
\rho_{i}=\max _{j=1, \ldots, J} \sup _{\lambda \in\left[b_{i}, B_{i}\right]}\left|g_{j}(\lambda)\right| .
$$

Since $\zeta_{r s}$ are i.i.d. $N(0,1)$, (55) implies that $\sum_{r=1}^{p} E\left|X_{p r}\right|^{2+\delta} \rightarrow 0$ as $p \rightarrow \infty$, which means that the Lyapunov condition holds for $X_{p r}$. As is well known, Lyapunov's condition implies Lindeberg's condition. Hence, condition a) of Lemma 14 is satisfied for $X_{p r}$.

Let us consider $\sum_{r=1}^{p} X_{p r}^{2}$. Since the convergence in mean implies the convergence in probability, the conditional Lindeberg condition is satisfied for $X_{p r}$ because the unconditional Lindeberg condition is satisfied as checked above. Further, in notations of Lemma 15, it is easy to see that

$$
V_{p p}^{2}=\sum_{j, j_{1}}\left[\left(\sum_{1 \leq s \leq t \leq q} \gamma_{j s t} \gamma_{j_{1} s t}\left(1+\delta_{s t}\right)\right) \frac{1}{p} \sum_{r=1}^{p} f_{j}\left(\mu_{r, i}\right) f_{j_{1}}\left(\mu_{r, i}\right)\right] .
$$

The convergence of the empirical distribution of $\mu_{1, i}, \ldots, \mu_{p, i}$ to $F\left(\lambda ; x_{i}\right)$ and the equality of $g_{j}$ and $f_{j}$ on the support of $F\left(\lambda ; x_{i}\right)$ implies that

$$
V_{p p}^{2} \xrightarrow{\mathrm{P}} \Sigma \equiv \sum_{j, j_{1}}\left[\left(\sum_{1 \leq s \leq t \leq q} \gamma_{j s t} \gamma_{j_{1} s t}\left(1+\delta_{s t}\right)\right) \int g_{j}(\lambda) g_{j_{1}}(\lambda) \mathrm{d} F\left(\lambda ; x_{i}\right)\right] .
$$

In particular, $V_{p p}^{2}$ is tight and Lemma 15 applies. Therefore, $\sum_{r=1}^{p} X_{p r}^{2}$ converges to the same limit as $V_{p p}^{2}$. Thus, by Lemma 14 , we get $\sum_{r=1}^{p} X_{p r} \xrightarrow{d} N(0, \Sigma)$.

Finally, let

$$
Y_{p r}=\frac{1}{\sqrt{p}} \sum_{(j, s, t) \in \Theta} \gamma_{j s t} g_{j}\left(\mu_{r, i}\right)\left(\zeta_{r s} \zeta_{r t}-\delta_{s t}\right)
$$

Since

$$
\operatorname{Pr}\left(\sum_{r=1}^{p} X_{p r} \neq \sum_{r=1}^{p} Y_{p r}\right) \rightarrow 0
$$

as $p \rightarrow \infty$, we have $\sum_{r=1}^{p} Y_{p r} \xrightarrow{d} N(0, \Sigma)$. Lemma DJO7 follows from this convergence via the Cramer-Wold device.
4.4. Derivation of equations (DJO27-29). Expression (DJO27) immediately follows from (DJO15). For (DJO28), differentiating identity (DJO13) with respect to $z$, we obtain

$$
1+\frac{\gamma_{1} m^{\prime}(z ; x)}{\left(1+\gamma_{1} m(z ; x)\right)^{2}}=\frac{m^{\prime}(z ; x)}{m^{2}(z ; x)}+\frac{-x^{2} \gamma_{2} m^{\prime}(z ; x)}{\left(1-\gamma_{2} x m(z ; x)\right)^{2}} .
$$

Setting $z=0$ and $x=x_{i}$, and using the fact that

$$
\begin{equation*}
m\left(0 ; x_{i}\right)=-\left(s_{i}+\gamma_{1}\right)^{-1}, \tag{56}
\end{equation*}
$$

which follows from (DJO15), we obtain

$$
1+\frac{\gamma_{1} m^{\prime}\left(0 ; x_{i}\right)}{\left(1-\gamma_{1}\left(s_{i}+\gamma_{1}\right)^{-1}\right)^{2}}=\frac{m^{\prime}\left(0 ; x_{i}\right)}{\left(s_{i}+\gamma_{1}\right)^{-2}}+\frac{-x_{i}^{2} \gamma_{2} m^{\prime}\left(0 ; x_{i}\right)}{\left(1+\gamma_{2} x_{i}\left(s_{i}+\gamma_{1}\right)^{-1}\right)^{2}} .
$$

Using the definition (DJO17) of $x_{i}$, we obtain

$$
1+\frac{\gamma_{1} m^{\prime}\left(0 ; x_{i}\right)}{\left(1-\gamma_{1}\left(s_{i}+\gamma_{1}\right)^{-1}\right)^{2}}=\frac{m^{\prime}\left(0 ; x_{i}\right)}{\left(s_{i}+\gamma_{1}\right)^{-2}}-\frac{\left(s_{i}+\gamma_{1}\right)^{2}\left(s_{i}+1\right)^{2} \gamma_{2} m^{\prime}\left(0 ; x_{i}\right)}{s_{i}^{2}}
$$

which implies (DJO28).
Finally, differentiating identity (DJO13) with respect to $x$, we obtain

$$
\frac{\gamma_{1} \mathrm{~d} m(z ; x) / \mathrm{d} x}{\left(1+\gamma_{1} m(z ; x)\right)^{2}}=\frac{\mathrm{d} m(z ; x) / \mathrm{d} x}{(m(z ; x))^{2}}+\frac{-1+\gamma_{2} x m(z ; x)-x\left(\gamma_{2} m(z ; x)+\gamma_{2} x \mathrm{~d} m(z ; x) / \mathrm{d} x\right)}{\left(1-\gamma_{2} x m(z ; x)\right)^{2}} .
$$

Setting $z=0$ and $x=x_{i}$, we obtain

$$
\frac{\gamma_{1} \mathrm{~d} m\left(0 ; x_{i}\right) / \mathrm{d} x}{\left(1+\gamma_{1} m\left(0 ; x_{i}\right)\right)^{2}}=\frac{\mathrm{d} m\left(0 ; x_{i}\right) / \mathrm{d} x}{\left(m\left(0 ; x_{i}\right)\right)^{2}}+\frac{-1-\gamma_{2} x_{i}^{2} \mathrm{~d} m\left(0 ; x_{i}\right) / \mathrm{d} x}{\left(1-\gamma_{2} x_{i} m\left(0 ; x_{i}\right)\right)^{2}} .
$$

This equality, the definition (DJO17) of $x_{i}$, and equation (56) imply (DJO29).

## 5. Analysis of the joint density of eigenvalues.

5.1. Proof of Lemma DJO13 about branch determination on the horizontal part of $\mathcal{K}$. To determine the branches, we will view the part of $\mathcal{K}$ on the real axis as the limit of a wedge-like contour

$$
\mathcal{W}=\left(\tilde{x}_{0}+\mathrm{i} \varepsilon, \tilde{\lambda}_{p 1}\right) \cup\left(\tilde{x}_{0}-\mathrm{i} \varepsilon, \tilde{\lambda}_{p 1}\right)
$$

as $\varepsilon \downarrow 0$, where i is the imaginary unit. Contour $\mathcal{W}$ intersects with each of the rays $\{z: \arg z=\varphi\}$, $\varphi \in(-\pi, \pi]$ no more than once, and therefore, the branches of all the fractional powers in $\omega$ must be principal as discussed in DJO. As $\varepsilon \downarrow 0$, we identify the branches by continuity as follows.
The situation will depend on which of $z_{1}, \ldots, z_{r}$ belong to the "upper" and which of them belong to the "lower" parts of $\mathcal{K} \cap\left[\tilde{x}_{0}, \tilde{\lambda}_{p 1}\right]$, that is the parts that are oriented from $\tilde{\lambda}_{p 1}$ to $\tilde{x}_{0}$, and from $\tilde{x}_{0}$ to $\tilde{\lambda}_{p 1}$, respectively.

There are $2^{r}$ possible scenarios: $\left(s g n_{z_{1}}= \pm 1, \ldots, s g n_{z_{r}}= \pm 1\right)$, where $s g n_{z_{j}}=+1$ means that $z_{j}$ belongs to the "upper" part, and $\operatorname{sgn} n_{j}=-1$ means that $z_{j}$ belongs to the "lower" part of $\mathcal{K} \cap\left[\tilde{x}_{0}, \tilde{\lambda}_{p 1}\right]$. Consider a particular scenario $\left(\operatorname{sgn} n_{z_{1}}, . ., \operatorname{sgn} n_{z_{r}}\right)$. Deforming $\mathcal{K} \cap\left[\tilde{x}_{0}, \tilde{\lambda}_{p 1}\right]$ to the wedgelike contour $\mathcal{W}$, we move $z_{j}$ to

$$
z_{j \varepsilon}=z_{j}+\mathrm{i} \times \operatorname{sgn}_{z_{j}} \frac{\tilde{\lambda}_{p 1}-z_{j}}{\tilde{\lambda}_{p 1}-\tilde{x}_{0}} \varepsilon .
$$

Since on $\mathcal{W}$, the principal branches of fractional powers are taken, the sign of the imaginary part of $\left(1-z_{j \varepsilon} z_{i \varepsilon}^{-1}\right)^{1 / 2}$ for $j>i$ must be equal to $\operatorname{sgn}_{z_{i}}$. Therefore, for $j>i$,

$$
\operatorname{sgn} \operatorname{Im}\left(1-z_{j} z_{i}^{-1}\right)^{1 / 2}=\lim _{\varepsilon \downarrow 0} \operatorname{sgn} \operatorname{Im}\left(1-z_{j \varepsilon} z_{i \varepsilon}^{-1}\right)^{1 / 2}=\operatorname{sgn}_{z_{i}} .
$$

Similarly, for $\tilde{\lambda}_{p s}$ and $z_{i}$ such that $\tilde{\lambda}_{p s}>z_{i}$, we have

$$
\operatorname{sgn} \operatorname{Im}\left(1-\tilde{\lambda}_{s} z_{i}^{-1}\right)^{-1 / 2}=\lim _{\varepsilon \downarrow 0} \operatorname{sgn} \operatorname{Im}\left(1-\tilde{\lambda}_{p s} z_{i \varepsilon}^{-1}\right)^{-1 / 2}=-\operatorname{sgn}_{z_{i}} .
$$

5.2. Proof of the fact that $\mathcal{I}_{\sigma}^{\prime}$ vanishes if there are any repeats in $\sigma$. Suppose, specifically, that $\sigma_{a+1}=\ldots=\sigma_{a+b}=i_{0} \leq k$ for some $b \geq 2$, so that the variables $z_{a+1}, \ldots, z_{a+b}$ lie within the same segment. Because of the branch effects described in Lemma DJO13, it helps the bookkeeping to first factor out from $\omega$ the terms that depend on $\operatorname{sgn}_{z_{i}}, i=1, \ldots, r$. To this end, let $v_{i}=\#\left\{s: \tilde{\lambda}_{p s}>z_{i}\right\}$ and $v_{+}=\sum_{i=1}^{r} v_{i}$. Then we factorize $\omega=S \tilde{\omega}$ so that $\tilde{\omega}$ does not depend on $\left\{\operatorname{sgn}_{z_{i}}, i=1, \ldots, r\right\}$ while from Lemma DJO13 we have

$$
S=\mathrm{i}^{r(r-1) / 2}(-\mathrm{i})^{v_{+}} \prod_{i=1}^{r} \operatorname{sgn}_{z_{i}}^{r-i+v_{i}} .
$$

It will be enough to show that the inner integral within (DJO43) given by

$$
\begin{equation*}
\mathcal{I}_{\sigma}^{\prime, \text { in }}=\int_{\mathcal{K}_{i_{0}}^{b} \cap \mathcal{S}^{b}} S \tilde{\omega} \mathcal{F}_{p, n} \prod_{i=a+1}^{a+b} \mathrm{~d} z_{i}=0 \tag{57}
\end{equation*}
$$

where $\mathcal{S}^{b}$ denotes the region $z_{a+1}<\ldots<z_{a+b}$. To this end, we factorize $S=S_{0} S_{1}$, where $S_{0}$ has all the terms that do not involve $z_{a+1}, \ldots, z_{a+b}$ and, if we let $v$ denote the common value of $v_{a+1}=\ldots=v_{a+b}$, we have

$$
S_{1}=\prod_{i=a+1}^{a+b} \operatorname{sgn}_{z_{i}}^{r-i+v_{i}} .
$$

Write $[\alpha, \beta]$ for the interval $\left[\tilde{\lambda}_{p, i_{0}+1}, \tilde{\lambda}_{p, i_{0}}\right]$. Recall that each of the $b$ copies of $\mathcal{K}_{i_{0}}$ is the union of two contour segments, namely $\mathcal{K}_{i_{0}}^{-}$, which traverses $[\alpha, \beta]$ left-to-right, and on which $\operatorname{sgn}_{z}=-1$, and $\mathcal{K}_{i_{0}}^{+}$, which traverses $[\alpha, \beta]$ right-to-left, and on which $s g n_{z}=1$. Let $s g n_{j}=s g n_{z_{a+j}}$ and decompose the inner integral (57) over all combinations of these contour segments to get

$$
\mathcal{I}_{\sigma}^{\prime, i n}=S_{0} \sum_{s g n_{1}} \cdots \sum_{s g n_{b}} \int_{[\alpha, \beta]^{b} \cap S^{b}} S_{1} S_{2} \tilde{\omega} \mathcal{F}_{p, n} \prod_{i=a+1}^{a+b} \mathrm{~d} z_{i},
$$

where each sum is over $s g n_{j}= \pm 1$ and the term $S_{2}=\prod_{j=1}^{b}\left(-s g n_{j}\right)$ counts whether $z_{a+j}$ traverses $\mathcal{K}_{i_{0}}^{-}$or $\mathcal{K}_{i_{0}}^{+}$. Now we can evaluate

$$
\sum_{s g n_{1}} \cdots \sum_{s g n_{b}} S_{1} S_{2}=(-1)^{b} \prod_{j=1}^{b} \sum_{s n_{j}} s g n_{j}^{r-a-j+v+1}=0
$$

when $b \geq 2$. Indeed, each sum on the right is of the form $1+(-1)^{l}$, which vanishes when $l$ is odd, and this must occur for at least one term if $b \geq 2$.
5.3. Proof of Lemma DJO16 about the Laplace approximation for $\mathcal{F}_{p, n}^{\operatorname{SigD}}$. Let $A=\operatorname{diag}\left\{a_{1}, \ldots, a_{k}\right\}$ with $0<a_{1}<\ldots<a_{k}$ and $B=\operatorname{diag}\left\{b_{1}, \ldots, b_{k}\right\}$ with $b_{1}>\ldots>b_{k}>0$. We have the following lemma.

Lemma 16. Let $c_{i j}=\left(a_{j}-a_{i}\right)\left(b_{i}-b_{j}\right) /\left\{\left(1+a_{i} b_{i}\right)\left(1+a_{j} b_{j}\right)\right\}$. Then, as $T \rightarrow \infty$,

$$
\begin{equation*}
{ }_{1} F_{0}(T / 2 ;-A, B)=\Gamma_{k}(k / 2) \pi^{-k(k+1) / 4} \prod_{i=1}^{k}\left(1+a_{i} b_{i}\right)^{-T / 2} \prod_{i<j}^{k}\left(T c_{i j} / 2\right)^{-1 / 2}(1+o(1)), \tag{58}
\end{equation*}
$$

where $o(1) \rightarrow 0$ uniforms on any compact subsects of the simplexes $0<a_{1}<\ldots<a_{k}$ and $b_{1}>\ldots>$ $b_{k}>0$.

This lemma is a minor extension of Chang's (1970) Theorem 1, which establishes (58) for fixed $A$ and $B$ (with both sides of (58) divided by the volume of the orthogonal group $\mathcal{O}(k)$ ). To show that the $o(1)$ is uniform on the set of $A$ and $B$ described in the lemma, it is sufficient to replace Hsu's (1948) Lemma 1 by Glynn's (1980) Theorem 2.1 in the proof of Chang (1970). Lemma DJO16 is a corollary to Lemma 16.
5.4. Proof of Lemma DJO17 about the Laplace approximation for $\mathcal{F}_{p, n}^{\mathrm{REG}}$. The proof below uses many ideas from Glynn (1980). First, let us represent $\mathcal{F}_{p, n}^{\mathrm{REG}}$ in terms of ${ }_{0} F_{1}$. Using the identities (see James' (1964) equations (30-31) or Glynn's equations (5.1-5.2))

$$
\mathcal{F}_{p, n}^{\mathrm{REG}}=\int_{\mathcal{O}(k)}{ }_{1} F_{1}\left(\alpha, \beta ; \frac{T}{2} s^{1 / 2} H Z H^{\prime} s^{1 / 2}\right)(\mathrm{d} H)
$$

where $\alpha=T a+\frac{k+1}{2}, \beta=T b+\frac{k+1}{2}$ and $(\mathrm{d} H)$ is the normalized invariant measure on the orthogonal group $\mathcal{O}(k)$, and

$$
{ }_{1} F_{1}\left(\alpha, \beta ; \frac{T}{2} s^{1 / 2} H Z H^{\prime} s^{1 / 2}\right)=\left(\Gamma_{k}(\alpha)\right)^{-1} \int_{\Sigma>0} \operatorname{etr}(-\Sigma)|\Sigma|^{T a}{ }_{0} F_{1}\left(\beta ; \frac{T}{2} s^{1 / 2} H Z H^{\prime} s^{1 / 2} \Sigma\right)(\mathrm{d} \Sigma),
$$

where $\Sigma>0$ is a positive deifinite $k \times k$ matrix, we obtain

$$
\mathcal{F}_{p, n}^{\mathrm{REG}}=\left(\Gamma_{k}(\alpha)\right)^{-1} \int_{\mathcal{O}(k)} \int_{\Sigma>0} \operatorname{etr}(-\Sigma)|\Sigma|^{T a}{ }_{0} F_{1}\left(\beta ; \frac{T}{2} s^{1 / 2} H Z H^{\prime} s^{1 / 2} \Sigma\right)(\mathrm{d} \Sigma)(\mathrm{d} H)
$$

Next, let us change variables of integration $\Sigma \mapsto Q, R$, where $\Sigma=\frac{T}{2} Q^{\prime} R^{2} Q$,

$$
R=\operatorname{diag}\left(r_{1}, \ldots, r_{k}\right)
$$

with $r_{1} \geq \ldots \geq r_{k} \geq 0$, and $Q \in \mathcal{O}(k)$. For this transformation, we have (see, for example, Herz (1955), p. 479)

$$
(\mathrm{d} \Sigma)=\mathcal{V}_{k}\left(\frac{T}{2}\right)^{\frac{k^{2}+k}{2}} 2^{k}|R| \prod_{i<j}\left(r_{i}^{2}-r_{j}^{2}\right)(\mathrm{d} R)(\mathrm{d} Q),
$$

where

$$
\mathcal{V}_{k}=2^{k} \pi^{k^{2} / 2} / \Gamma_{k}(k / 2)
$$

is the volume of $\mathcal{O}(k)$, and $(\mathrm{d} Q)$ is the normalized invariant measure on $\mathcal{O}(k)$.

The transformation is one-to- $2^{k}$ becase $Q$ is only determined up to a left-multiplication by a diagonal matrix with $\pm 1$ coefficients along the diagonal. Therefore, we have

$$
\begin{align*}
\mathcal{F}_{p, n}^{\mathrm{REG}}= & C_{T} \int_{\mathcal{O}(k) \times D(R) \times \mathcal{O}(k)} \operatorname{etr}\left(-T R^{2} / 2\right)|R|^{2 \alpha-k} \prod_{i<j}\left(r_{i}^{2}-r_{j}^{2}\right) \\
& \times{ }_{0} F_{1}\left(\beta ; \frac{T^{2}}{4} R Q s^{1 / 2} H Z H^{\prime} s^{1 / 2} Q^{\prime} R\right)(\mathrm{d} Q)(\mathrm{d} R)(\mathrm{d} H) \tag{59}
\end{align*}
$$

where

$$
C_{T}=\mathcal{V}_{k}(T / 2)^{k \alpha} / \Gamma_{k}(\alpha), \text { and } D(R)=\left\{R: r_{1}>\ldots>r_{k}>0\right\} .
$$

Consider Herz' integral representation for ${ }_{0} F_{1}$ (see Butler and Wood (2003), equation (12))

$$
\begin{equation*}
{ }_{0} F_{1}\left(\beta, \frac{1}{4} \Theta \Theta^{\prime}\right)=c \int_{\mathcal{U}} \operatorname{etr}\{\Theta Y\}\left|I_{k}-Y Y^{\prime}\right|^{(2 \beta-2 k-1) / 2}(\mathrm{~d} Y), \tag{60}
\end{equation*}
$$

where $\Theta$ is a $k \times k$ matrix, $\mathcal{U}=\left\{Y: Y Y^{\prime}<I_{k}\right\}$ and $c=\pi^{-k^{2} / 2} \Gamma_{k}(\beta) / \Gamma_{k}(\beta-k / 2)$. This representation is valid when $\beta \geq k$, which holds for sufficiently large $T$. Using (60) in (59), we obtain

$$
\begin{aligned}
\mathcal{F}_{p, n}^{\mathrm{REG}}= & \tilde{C}_{T} \int_{\mathcal{O}(k) \times D(R) \times \mathcal{O}(k) \times \mathcal{U}} \operatorname{etr}\left[-T R^{2} / 2+T R Q s^{1 / 2} H Z^{1 / 2} Y\right] \\
& \times|R|^{2 T a+1} \prod_{i<j}\left(r_{i}^{2}-r_{j}^{2}\right)\left|I_{k}-Y Y^{\prime}\right|^{T b-k / 2} \mathrm{~d} z,
\end{aligned}
$$

where

$$
\mathrm{d} z=(\mathrm{d} Q)(\mathrm{d} R)(\mathrm{d} H)(\mathrm{d} Y), \text { and } \tilde{C}_{T}=\frac{2^{k}(T / 2)^{k \alpha} \Gamma_{k}(\beta)}{\Gamma_{k}(\alpha) \Gamma_{k}(k / 2) \Gamma_{k}(\beta-k / 2)} .
$$

Now let us make the change of variables $Y \mapsto U, V, \sigma$, where $Y=U \sigma V$ is a singular value decomposition of $Y$ with

$$
\sigma=\operatorname{diag}\left(\sigma_{1}, \ldots, \sigma_{k}\right), \text { with } \sigma_{1} \geq \ldots \geq \sigma_{k} \geq 0
$$

For such a change of variables, we have

$$
(\mathrm{d} Y)=\mathcal{V}_{k}^{2} \prod_{i<j}\left(\sigma_{i}^{2}-\sigma_{j}^{2}\right)(\mathrm{d} U)(\mathrm{d} V)(\mathrm{d} \sigma)
$$

where $(\mathrm{d} U)$ and $(\mathrm{d} V)$ are the normalized invariant measures on $\mathcal{O}_{k}$.
The transformation $Y \mapsto U, V, \sigma$ is one-to- $2^{k}$. Therefore, we obtain

$$
\begin{align*}
\mathcal{F}_{p, n}^{\mathrm{REG}}= & \hat{C}_{T} \int_{\Lambda} \operatorname{etr}\left[T\left(-R^{2} / 2+R Q s^{1 / 2} H Z^{1 / 2} U \sigma V\right)\right] \\
& \times|R|^{2 T a+1} \prod_{i<j}\left(r_{i}^{2}-r_{j}^{2}\right)\left(\sigma_{i}^{2}-\sigma_{j}^{2}\right)\left|I_{k}-\sigma^{2}\right|^{T b-k / 2} \mathrm{~d} \xi, \tag{61}
\end{align*}
$$

where

$$
\Lambda=\mathcal{O}(k) \times D(R) \times \mathcal{O}(k) \times \bar{D}(\sigma) \times \mathcal{O}(k) \times \mathcal{O}(k)
$$

with $\bar{D}(\sigma)=\left\{\sigma: 1 \geq \sigma_{1} \geq \ldots \geq \sigma_{k} \geq 0\right\}$,

$$
\mathrm{d} \xi=(\mathrm{d} Q)(\mathrm{d} R)(\mathrm{d} H)(\mathrm{d} \sigma)(\mathrm{d} U)(\mathrm{d} V), \text { and } \hat{C}_{T}=\frac{2^{2 k} \pi^{k^{2}}(T / 2)^{k \alpha} \Gamma_{k}(\beta)}{\Gamma_{k}(\alpha)\left(\Gamma_{k}(k / 2)\right)^{3} \Gamma_{k}(\beta-k / 2)} .
$$

Equation (61) can be rewritten in the form amenable to the Laplace approximation method as follows

$$
\begin{equation*}
\mathcal{F}_{p, n}^{\mathrm{REG}}=\hat{C}_{T} \int_{\Lambda} h f^{T} \mathrm{~d} \xi \tag{62}
\end{equation*}
$$

where

$$
f=\operatorname{etr}\left\{-R^{2} / 2+R Q s^{1 / 2} H Z^{1 / 2} U \sigma V\right\}\left|I_{k}-\sigma^{2}\right|^{b}|R|^{2 a}
$$

and

$$
h=\left|I_{k}-\sigma^{2}\right|^{-k / 2}|R| \prod_{i<j}\left(r_{i}^{2}-r_{j}^{2}\right)\left(\sigma_{i}^{2}-\sigma_{j}^{2}\right)
$$

By Lemma 4.2 of Glynn (1980), the maximum of $\operatorname{etr}\left\{R Q s^{1 / 2} H Z^{1 / 2} U \sigma V\right\}$ over $(Q, H, U, V) \in$ $\mathcal{O}(k)^{4}$ is achieved at $2^{3 k}$ points, where $Q, H, U$, and $V$ are diagonal with values $\pm 1$ along the diagonal, and such that $Q H U V=I_{k}$. The value of $\operatorname{etr}\left\{R Q s^{1 / 2} H Z^{1 / 2} U \Sigma V\right\}$ at the maximum is $\operatorname{etr}\left\{\sum_{j=1}^{k} r_{j} \sigma_{j} \sqrt{2 \zeta_{j}}\right\}$, where

$$
\zeta_{j}=z_{j} s_{j} / 2
$$

We introduce this notation because it simplifies some expressions later on.
This implies that the maximum of $f$ over $\Lambda$ is achieved at $2^{3 k}$ points with $Q, H, U, V$ as above, and with $r_{j}$ and $\sigma_{j}$ satisfying the following first order conditions for maximisation with respect to $r_{j}$ and $\sigma_{j}, j=1, \ldots, k$,

$$
\begin{equation*}
-r_{j}+\sigma_{j} \sqrt{2 \zeta_{j}}+2 a / r_{j}=0 \tag{63}
\end{equation*}
$$

and

$$
\begin{equation*}
r_{j} \sqrt{2 \zeta_{j}}-\frac{2 b \sigma_{j}}{1-\sigma_{j}^{2}}=0 \tag{64}
\end{equation*}
$$

From (63) we obtain

$$
\sigma_{j}=\frac{r_{j}^{2}-2 a}{r_{j} \sqrt{2 \zeta_{j}}}
$$

Using this in (64) we obatin

$$
\left(r_{j}^{2}-2 a\right)^{2}+2\left(b-\zeta_{j}\right)\left(r_{j}^{2}-2 a\right)-4 a \zeta_{j}=0
$$

Let

$$
z_{j+}=\frac{1}{2}\left\{\zeta_{j}-b+\sqrt{\left(b-\zeta_{j}\right)^{2}+4 a \zeta_{j}}\right\}
$$

be the positive solution of the quadratic equation $z_{j}^{2}+\left(b-\zeta_{j}\right) z_{j}-a \zeta_{j}=0$. Then, $r_{j}^{2}-2 a=2 z_{j+}$, and thus

$$
\begin{align*}
r_{j}^{2} & =2\left(z_{j+}+a\right), \text { and }  \tag{65}\\
\sigma_{j}^{2} & =\frac{1}{\zeta_{j}} \frac{z_{j+}^{2}}{z_{j+}+a} \tag{66}
\end{align*}
$$

We have introduced the notation $z_{j+}$ here for compatibility with the previous version of DJO, that studied the one dimensional case $k=1$.

Let us verify the second order conditions for the maximum at the above $r_{j}^{2}$ and $\sigma_{j}^{2}$. The matrix of the second derivatives of $\ln f$ with respect to $r_{j}$ and $\sigma_{j}$ is

$$
\left(\begin{array}{cc}
-1-2 a / r_{j}^{2} & \sqrt{2 \zeta_{j}} \\
\sqrt{2 \zeta_{j}} & -2 b\left(1+\sigma_{j}^{2}\right) /\left(1-\sigma_{j}^{2}\right)^{2}
\end{array}\right) .
$$

Its value at the critical point is

$$
\left(\begin{array}{cc}
-\left(z_{j+}+2 a\right) /\left(z_{j+}+a\right) & \sqrt{2 \zeta_{j}} \\
\sqrt{2 \zeta_{j}} & -2 \zeta_{j}\left(z_{j+}+a\right) / b-2 \zeta_{j}^{2}\left(z_{j+}+a\right)^{2} /\left(b z_{j+}^{2}\right)
\end{array}\right) .
$$

By inspection, the diagonal elements of this matrix are negative, whereas the determinant

$$
4 \zeta_{j}\left(a+z_{j+}\right)\left(a \zeta_{j}+z_{j+}^{2}\right) /\left(b z_{j+}^{2}\right)>0
$$

so that the second order condition for the maximum is satisfied.
Note that the value of $\ln f$ at the maximum is

$$
\max \ln f=\sum_{j=1}^{k}\left\{-r_{j}^{2} / 2+r_{j} \sigma_{j} \sqrt{2 \zeta_{j}}+b \ln \left(1-\sigma_{j}^{2}\right)+2 a \ln r_{j}\right\},
$$

where $r_{j}>0$ and $\sigma_{j}>0$ are given by (65) and (66). Expressing max $\ln f$ in terms of $a, b, \zeta_{j}$ and $z_{j+}$, we obtain

$$
\max \ln f=\sum_{j=1}^{k}\left\{z_{j+}-b \ln \left(z_{j+}+b\right)+a \ln \left(z_{j+}+a\right)+a \ln 2+b \ln b-a\right\}
$$

Since the maximum is achieved at the $2^{3 k}$ points, the integral over $\Lambda$ in (62) can be replaced, for the purpose of the asymptotic analysis, by $2^{3 k}$ times the integral over

$$
\Lambda^{+}=\mathcal{O}^{+}(k) \times D(R) \times \mathcal{O}^{+}(k) \times \bar{D}(\sigma) \times \mathcal{O}^{+}(k) \times \mathcal{O}^{+}(k),
$$

where $\mathcal{O}^{+}(k)$ denotes the set of $k$-dimensional orthogonal matrices with positive diagonal elements. Since $Q, H, U$, and $V$ are proper in $\mathcal{O}^{+}(k)$, they can be parameterized as

$$
Q=\exp \{G\}, H=\exp \{F\}, U=\exp \{A\}, \text { and } V=\exp \{B\},
$$

where $G, F, A$, and $B$ are $k \times k$ skew symmetric.
Anderson (1965) shows that the Jacobian of the transfomation $Q \mapsto S$ equals

$$
J_{1}=J(Q \mapsto G)=\Gamma_{k}(k / 2) 2^{-k} \pi^{-k^{2} / 2}\left(1+O\left(G_{i j}^{2}\right)\right),
$$

where $O\left(G_{i j}^{2}\right)$ denotes the terms that are at least quadratic in the elements of $G$. Similar expressions hold for the Jacobians $J_{2}, J_{3}$ and $J_{4}$ of the transformations $H \mapsto F, U \mapsto A$, and $V \mapsto B$, respectively. By making this change of variables, we arrive at the following asymptotic representation

$$
\begin{equation*}
\mathcal{F}_{p, n}^{\mathrm{REG}} \sim c_{T} \int_{\Xi} g f^{T}(\mathrm{~d} \tilde{\xi}), \tag{67}
\end{equation*}
$$

where

$$
c_{n}=2^{3 k} \hat{C}_{T}=\frac{2^{5 k} \pi^{k^{2}}(T / 2)^{k \alpha} \Gamma_{k}(\beta)}{\Gamma_{k}(\alpha)\left(\Gamma_{k}(k / 2)\right)^{3} \Gamma_{k}(\beta-k / 2)},
$$

$\Xi$ is the image of $\Lambda^{+}$under the transformation $Q, H, U, V \mapsto G, F, A, B$,

$$
\begin{aligned}
g & =J_{1} J_{2} J_{3} J_{4}\left|I_{k}-\sigma^{2}\right|^{-k / 2}|R| \prod_{i<j}\left(r_{i}^{2}-r_{j}^{2}\right)\left(\sigma_{i}^{2}-\sigma_{j}^{2}\right), \text { and } \\
f & =\operatorname{etr}\left\{-R^{2} / 2+R e^{G} s^{1 / 2} e^{F} Z^{1 / 2} e^{A} \sigma e^{B}\right\}\left|I_{k}-\sigma^{2}\right|^{b}|R|^{2 a} .
\end{aligned}
$$

Expanding $e^{G}, e^{F}, e^{A}$, and $e^{B}$ into powers of $G, F, A$, and $B$, we have $\operatorname{tr}\left[R e^{G} s^{1 / 2} e^{F} Z^{1 / 2} e^{A} \sigma e^{B}\right]$ equals

$$
\begin{aligned}
& \sum_{j=1}^{k} r_{j} s_{j}^{1 / 2} z_{j}^{1 / 2} \sigma_{j}-\frac{1}{2} \sum_{j<i}^{k}\left(r_{j} s_{j}^{1 / 2} z_{j}^{1 / 2} \sigma_{j}+r_{i} s_{i}^{1 / 2} z_{i}^{1 / 2} \sigma_{i}\right)\left(G_{j i}^{2}+F_{j i}^{2}+A_{j i}^{2}+B_{j i}^{2}\right) \\
& -\sum_{j<i}^{k}\left(r_{j} s_{i}^{1 / 2} z_{j}^{1 / 2} \sigma_{j}+r_{i} s_{j}^{1 / 2} z_{i}^{1 / 2} \sigma_{i}\right) G_{j i} F_{j i}-\sum_{j<i}^{k}\left(r_{j} s_{i}^{1 / 2} z_{i}^{1 / 2} \sigma_{j}+r_{i} s_{j}^{1 / 2} z_{j}^{1 / 2} \sigma_{i}\right) G_{j i} A_{j i} \\
& -\sum_{j<i}^{k}\left(r_{j} s_{i}^{1 / 2} z_{i}^{1 / 2} \sigma_{i}+r_{i} s_{j}^{1 / 2} z_{j}^{1 / 2} \sigma_{j}\right) G_{j i} B_{j i}-\sum_{j<i}^{k}\left(r_{j} s_{j}^{1 / 2} z_{i}^{1 / 2} \sigma_{j}+r_{i} s_{i}^{1 / 2} z_{j}^{1 / 2} \sigma_{i}\right) F_{j i} A_{j i} \\
& -\sum_{j<i}^{k}\left(r_{j} s_{j}^{1 / 2} z_{i}^{1 / 2} \sigma_{i}+r_{i} s_{i}^{1 / 2} z_{j}^{1 / 2} \sigma_{j}\right) F_{j i} B_{j i}-\sum_{j<i}^{k}\left(r_{j} s_{j}^{1 / 2} z_{j}^{1 / 2} \sigma_{i}+r_{i} s_{i}^{1 / 2} z_{i}^{1 / 2} \sigma_{j}\right) A_{j i} B_{j i} \\
& \text { +h.o.t., }
\end{aligned}
$$

where h.o.t. stands for higher order terms (in $G_{j i}, F_{j i}, A_{j i}$, and $B_{j i}$ ).
The maximum value of $f$ is obtained at a single point $\tilde{\xi}$ in the interior of $\Xi$. The Hessian of $-\ln f$, $\Delta$, reduces to a product of determinants of matrices which are at most $4 \times 4$. A direct calculation (using MAPLE symbolic algebra) gives

$$
\begin{aligned}
\Delta= & \prod_{j=1}^{k} \frac{4 \zeta_{j}\left(a+z_{j+}\right)\left(a \zeta_{j}+z_{j+}^{2}\right)}{b z_{j+}^{2}} \prod_{j<i}^{k}\left\{\left(z_{i}-z_{j}\right)\left(s_{i}-s_{j}\right)\left(r_{i}^{2}-r_{j}^{2}\right)\right. \\
& \left.\left(\sigma_{i}^{2}-\sigma_{j}^{2}\right)\left(s_{i}^{1 / 2} z_{i}^{1 / 2} \sigma_{i} r_{i}-s_{j}^{1 / 2} z_{j}^{1 / 2} \sigma_{j} r_{j}\right)^{2}\right\} .
\end{aligned}
$$

Further, using Stirling's formula, we have, at $G=F=A=B=0$,

$$
c_{T} J_{1} J_{2} J_{3} J_{4} \sim \frac{\Gamma_{k}(k / 2)}{\pi^{k^{2}+k(k+1) / 4}} \frac{T^{k(3 k+1) / 4}}{2^{k T a+k^{2} / 2}} \frac{b^{k^{2} / 2} e^{k T a}}{a^{k T a+k(k+1) / 4}} .
$$

The statement of Lemma DJO17 now follows by applying Glynn's (1980) Theorem 2.1 to (67). We obtain

$$
\begin{aligned}
\mathcal{F}_{p, n}^{\mathrm{REG}} \sim & \left(\frac{T}{2}\right)^{-k(k-1) / 4} \frac{\Gamma_{k}(k / 2)}{\pi^{k(k+1) / 4}} \frac{b^{k T b+k(k+1) / 4}}{a^{k T a+k(k+1) / 4}} \\
& \times \prod_{j=1}^{k} e^{T z_{j}+} \frac{\left(z_{j+}+a\right)^{a T}}{\left(z_{j+}+b\right)^{b T}}\left(\frac{z_{j+}\left(z_{j+}+a\right)}{z_{j+}^{2}+a \zeta_{j}}\right)^{\frac{1}{2}} \\
& \times \prod_{i<j}\left(z_{i}-z_{j}\right)^{-1 / 2}\left(s_{i}-s_{j}\right)^{-1 / 2}
\end{aligned}
$$

By Glynn's theorem, this asymptotic approximation is uniform over $a$ and $b$ that belong to compact subsets of $(1 / 2, \infty)$ and $(0,1 / 2)$, respectively, and over $s$ and $Z$, such that their diagonal entries belong to compact subsets of the simplexes $s_{1}>\ldots>s_{k}>0$ and $z_{1}>\ldots>z_{k}>0$.
5.5. Proof of Lemma DJO18 about the Laplace approximation for $\overline{\mathcal{I}}_{\sigma \varnothing}^{\mathrm{SigD}}$. As shown in DJO, the inner-most integral in the multiple integral representation (DJO49) of $\overline{\mathcal{I}}_{\sigma_{\varnothing}}^{\mathrm{SigD}}$ equals

$$
\begin{equation*}
\overline{\mathcal{I}}_{\sigma_{\varepsilon}, i n}^{\mathrm{SigD}}=G_{-1}^{\mathrm{SigD}} \int_{\tilde{\lambda}_{p 2}+\varepsilon}^{\tilde{\lambda}_{p 1}} e^{-p f_{1}\left(z_{1}\right)} g_{1}\left(z_{1}\right)(1+o(1)) \mathrm{d} z_{1} . \tag{68}
\end{equation*}
$$

First, we will apply the Laplace method (see Olver (1997), p. 81-82) to the integral

$$
I(p) \equiv \int_{\tilde{\lambda}_{p 2}+\varepsilon}^{\tilde{\lambda}_{p 1}} e^{-p f_{1}\left(z_{1}\right)} g_{1}\left(z_{1}\right) \mathrm{d} z_{1}
$$

To line up our analysis with that on pages 81-82 of Olver (1997) and to simplify notations, let us rewrite $I(p)$ as

$$
I(p)=\int_{-\tilde{\lambda}_{p 1}}^{-\tilde{\lambda}_{p 2}-\varepsilon} e^{-p f(t)} g(t) \mathrm{d} t
$$

where $t=-z_{1}$,

$$
\begin{aligned}
& f(t)=\frac{N-p+2}{2 p} \ln \left(1+\frac{s_{1} t}{1+s_{1}}\right)+\frac{1}{2 p} \sum_{j=k+1}^{p} \ln \left(-t-\tilde{\lambda}_{p j}\right), \text { and } \\
& g(t)=\left(\tilde{\lambda}_{p 1}+t\right)^{-1 / 2} \prod_{j=2}^{k}\left(\left(-t-z_{j}\right) /\left(-t-\tilde{\lambda}_{p j}\right)\right)^{1 / 2}
\end{aligned}
$$

Let us show that under the null hypothesis, $f(t)$ has positive continuous derivative for $t$ from the interval $\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right]$ for sufficiently large $p$, a.s. All statements made in this section should be understood as holding for sufficiently large $p$, almost surely, and we will omit this qualification to avoid frequent repetitions. We have

$$
\begin{equation*}
f^{\prime}(t)=\frac{N-p+2}{2 p} \frac{s_{1}}{1+s_{1}+s_{1} t}+\frac{1}{2 p} \sum_{j=k+1}^{p}\left(t+\tilde{\lambda}_{p j}\right)^{-1} \tag{69}
\end{equation*}
$$

Clearly, $f^{\prime}(t)$ is continuous and strictly decreasing for all $t \in\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right]$. Under the null, the minimum of $f^{\prime}(t)$ on $t \in\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right]$ converges to

$$
\begin{equation*}
\frac{\rho^{2} s_{01}}{2 \gamma_{1} \gamma_{2}\left(1+s_{01}-s_{01}\left(\lim \tilde{\lambda}_{p 2}+\varepsilon\right)\right)}+\frac{1}{2} \tilde{m}\left(\lim \tilde{\lambda}_{p 2}+\varepsilon\right) \tag{70}
\end{equation*}
$$

uniformly over $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, a.s. Here $\tilde{m}(x)$ is the Stieltjes transform (or rather its analytic continuation to a point on the real line) of the limiting spectral distribution of the multivariate Beta matrix $\left(n_{H} \mathbf{H}+n_{2} \mathbf{E}\right)^{-1} n_{H} \mathbf{H}$ and

$$
\lim \tilde{\lambda}_{p 2} \equiv \frac{\left(s_{02}+\gamma_{1}\right)\left(s_{02}+1\right) \gamma_{2}}{s_{02}\left(\gamma_{1}+s_{02} \gamma_{2}+\gamma_{2}\right)}
$$

is the a.s. limit of $\tilde{\lambda}_{p 2}$ (see equation (DJO41)).

Consider the following function of two real variables

$$
\Psi\left(y, s_{01}\right) \equiv \frac{\rho^{2} s_{01}}{2 \gamma_{1} \gamma_{2}\left(1+s_{01}-s_{01} y\right)}+\frac{1}{2} \tilde{m}(y)
$$

From the above discussion, we see that this is the value of the a.s. limit of $f^{\prime}(t)$ as $t \rightarrow-y$. The function is well defined for all $y \leq 1$ that lie above the upper boundary of the support of the limiting spectral distribution of the multivariate Beta matrix, and for all $s_{01}>\bar{s}$. It is also well defined for positive $s_{01}$ and all, but one, values $y>1$, but we focus on the supercritical spikes (hence $s_{01}>\bar{s}$ ) and on $y$ that may be equal to a limit of an eigenvalue of a multivariate Beta matrix (hence $y \leq 1$ ).

Recall that $\tilde{\lambda}_{p j}=\alpha_{n} \lambda_{p j} /\left(1+\alpha_{n} \lambda_{p j}\right)$, where $\alpha_{n}=n_{H} / n_{2}$ and $\lambda_{p j}$ is the $j$-th largest eigenvalue of $\mathbf{F}=\mathbf{E}^{-1} \mathbf{H}$. This implies that the upper boundary of the support of the limiting spectral distribution of the multivariate Beta matrix $\left(n_{H} \mathbf{H}+n_{2} \mathbf{E}\right)^{-1} n_{H} \mathbf{H}$ equals

$$
\tilde{\beta}_{+} \equiv \alpha \beta_{+} /\left(1+\alpha \beta_{+}\right), \text {where } \alpha=\gamma_{2} / \gamma_{1}
$$

and that

$$
\begin{equation*}
\tilde{m}(y)=\alpha x+1+\alpha^{-1}(\alpha x+1)^{2} m(x), \tag{71}
\end{equation*}
$$

where $y=\alpha x /(1+\alpha x)$ and $m(x)$ is the Stieltjes transform of the limiting spectral distribution of F. It is well known (see, for example, p. 79 of Bai and Silverstein's (2010) book), that

$$
\begin{equation*}
m(x)=\frac{1}{x \gamma_{1}}-\frac{1}{x}-\frac{\gamma_{1}\left(x\left(1-\gamma_{2}\right)+1-\gamma_{1}\right)+2 x \gamma_{2}-\gamma_{1} \sqrt{\left(\left(1-\gamma_{1}\right)+x\left(1-\gamma_{2}\right)\right)^{2}-4 x}}{2 x \gamma_{1}\left(\gamma_{1}+x \gamma_{2}\right)} \tag{72}
\end{equation*}
$$

Now returning to function $\Psi$, note that this function is increasing in both $y$ and $s_{01}$ on $y \in\left(\tilde{\beta}_{+}, 1\right]$ and $s_{01}>\bar{s}$. Let us compute its limit as $y \rightarrow \tilde{\beta}_{+}$and $s_{01} \rightarrow \bar{s}$. Then $x \equiv y /(\alpha-\alpha y) \rightarrow \beta_{+} \equiv$ $\left((1+\rho) /\left(1-\gamma_{2}\right)\right)^{2}$, the square root in the above definition of $m$ converges to zero, and after some algebra, we obtain

$$
\frac{1}{2} \tilde{m}(y) \rightarrow-\frac{\rho\left((1+\rho)^{2} \gamma_{2}+\gamma_{1}\left(1-\gamma_{2}\right)^{2}\right)}{2 \gamma_{1} \gamma_{2}\left(1-\gamma_{2}\right)(1+\rho)}
$$

On the other hand, as $s_{01} \rightarrow \bar{s} \equiv\left(\gamma_{2}+\rho\right) /\left(1-\gamma_{2}\right)$ and $y \rightarrow \tilde{\beta}_{+}$,

$$
\frac{\rho^{2} s_{01}}{2 \gamma_{1} \gamma_{2}\left(1+s_{01}-s_{01} y\right)} \rightarrow \frac{\rho\left((1+\rho)^{2} \gamma_{2}+\gamma_{1}\left(1-\gamma_{2}\right)^{2}\right)}{2 \gamma_{1} \gamma_{2}\left(1-\gamma_{2}\right)(1+\rho)}
$$

so that $\Psi\left(y, s_{01}\right) \rightarrow 0$. Since $\Psi\left(y, s_{01}\right)$ is increasing, the limit $(70)$ of $f^{\prime}(t)$ must be positive. Moreover, there exists a small positive number such that $f^{\prime}(t)$ is above this number for all $t \in\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right]$ and all $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, for sufficiently large $p$, almost surely.

Since $f(t)$ is strictly increasing on $\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right]$, the main contribution to the integral $I(p)$ comes from the vicinity of $-\tilde{\lambda}_{p 1}$. From (69), we see that

$$
\begin{equation*}
f^{\prime}\left(-\tilde{\lambda}_{p 1}\right) \stackrel{\text { a.s. }}{ } \frac{\rho^{2} s_{01}}{2 \gamma_{1} \gamma_{2}\left(1+s_{01}-s_{01} \lim \tilde{\lambda}_{p 1}\right)}+\frac{1}{2} \tilde{m}\left(\lim \tilde{\lambda}_{p 1}\right), \tag{73}
\end{equation*}
$$

where

$$
\lim \tilde{\lambda}_{p 1} \equiv \frac{\left(s_{01}+\gamma_{1}\right)\left(s_{01}+1\right) \gamma_{2}}{s_{01}\left(\gamma_{1}+s_{01} \gamma_{2}+\gamma_{2}\right)}
$$

Using (71) and (72), we obtain, after some algebra

$$
\begin{equation*}
\tilde{m}\left(\lim \tilde{\lambda}_{p 1}\right)=\frac{\rho^{2}}{\gamma_{1} \gamma_{2}} \frac{s_{01}\left(\gamma_{1}+\gamma_{2}+\gamma_{2} s_{01}\right)}{\left(\gamma_{1}+s_{01}\right)\left(\gamma_{2}+\gamma_{2} s_{01}-s_{01}\right)} . \tag{74}
\end{equation*}
$$

Further,

$$
\frac{\rho^{2} s_{01}}{\gamma_{1} \gamma_{2}\left(1+s_{01}-s_{01} \lim \tilde{\lambda}_{p 1}\right)}=\frac{s_{01}\left(\gamma_{1}+\gamma_{2}+\gamma_{2} s_{01}\right)}{\gamma_{1} \gamma_{2}\left(1+s_{01}\right)} .
$$

These equations together with (73) yield

$$
\begin{equation*}
f^{\prime}\left(-\tilde{\lambda}_{p 1}\right) \xrightarrow{a . s .} H_{1} \equiv \frac{\left(1-\gamma_{2}\right) s_{01}\left(\gamma_{1}+\gamma_{2}+\gamma_{2} s_{01}\right)\left(1+s_{01}-\beta_{+}^{1 / 2}\right)\left(1+s_{01}-\beta_{-}^{1 / 2}\right)}{2 \gamma_{1} \gamma_{2}\left(\gamma_{1}+s_{01}\right)\left(1+s_{01}\right)\left(s_{01}-\gamma_{2} s_{01}-\gamma_{2}\right)} . \tag{75}
\end{equation*}
$$

The rest of our proof of Lemma DJO18 closely follows Olver (1997), p. 81-82. Consider a new variable of integration

$$
v=f(t)-f\left(-\tilde{\lambda}_{p 1}\right)
$$

and let $c$ be a fixed number that belongs to the a.s. limit of $\left(-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right)$ as $p, n \rightarrow_{\gamma} \infty$, under the null hypothesis. We have

$$
\begin{equation*}
e^{p f\left(-\tilde{\lambda}_{p 1}\right)} \int_{-\tilde{\lambda}_{p 1}}^{c} e^{-p f(t)} g(t) \mathrm{d} t=\int_{0}^{\kappa} e^{-p v} q(v) \mathrm{d} v \tag{76}
\end{equation*}
$$

where

$$
\kappa=f(c)-f\left(-\tilde{\lambda}_{p 1}\right) \text { and } q(v)=g(t)(\mathrm{d} t / \mathrm{d} v)=g(t) / f^{\prime}(t) .
$$

Note that there exist fixed non-random $0<K_{1}<K_{2}<\infty$, such that $\kappa \in\left[K_{1}, K_{2}\right]$.
By definition of $f(t)$ and $g(t)$, as $t \rightarrow-\tilde{\lambda}_{p 1}$,

$$
\left(f(t)-f\left(-\tilde{\lambda}_{p 1}\right)\right) /\left(t+\tilde{\lambda}_{p 1}\right) \rightarrow F, \text { and } g(t) /\left(t+\tilde{\lambda}_{p 1}\right)^{-1 / 2} \rightarrow G,
$$

where

$$
F \equiv f^{\prime}\left(-\tilde{\lambda}_{p 1}\right), \text { and } G \equiv \prod_{j=2}^{k}\left(\left(\tilde{\lambda}_{p 1}-z_{j}\right) /\left(\tilde{\lambda}_{p 1}-\tilde{\lambda}_{p j}\right)\right)^{1 / 2}
$$

These convergences, together with the above definitions of $v$ and $q(v)$ imply that

$$
\begin{equation*}
q(v) \sim G v^{-1 / 2} / F^{1 / 2} \text { as } v \rightarrow 0+ \tag{77}
\end{equation*}
$$

This means that the ratio of the left hand side of (77) to its right hand side converges to 1 as $v \rightarrow 0+$. This convergece is uniform in $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$.

Now, closely following Olver, rearrange the integral (76) in the form

$$
\begin{equation*}
\int_{0}^{\kappa} e^{-p v} q(v) \mathrm{d} v=G F^{-1 / 2}\left\{\int_{0}^{\infty} e^{-p v} v^{-1 / 2} \mathrm{~d} v-\varepsilon_{1}(p)\right\}+\varepsilon_{2}(p), \tag{78}
\end{equation*}
$$

where

$$
\varepsilon_{1}=\int_{\kappa}^{\infty} e^{-p v} v^{-1 / 2} \mathrm{~d} v \text { and } \varepsilon_{2}=\int_{0}^{\kappa} e^{-p v}\left\{q(v)-G F^{-1 / 2} v^{-1 / 2}\right\} \mathrm{d} v .
$$

Since $\kappa \in\left[K_{1}, K_{2}\right]$ for all sufficiently large $p$,

$$
\int_{0}^{\infty} e^{-p v} v^{-1 / 2} \mathrm{~d} v-\varepsilon_{1}(p) \sim \Gamma(1 / 2) p^{-1 / 2}=\sqrt{\pi / p} \text { as } p \rightarrow \infty .
$$

Further, the above results on the derivative of $f$ yield

$$
F \sim H_{1} \equiv \frac{\left(1-\gamma_{2}\right) s_{01}\left(\gamma_{1}+\gamma_{2}+\gamma_{2} s_{01}\right)\left(1+s_{01}-\beta_{+}^{1 / 2}\right)\left(1+s_{01}-\beta_{-}^{1 / 2}\right)}{2 \gamma_{1} \gamma_{2}\left(\gamma_{1}+s_{01}\right)\left(1+s_{01}\right)\left(s_{01}-\gamma_{2} s_{01}-\gamma_{2}\right)}
$$

Therefore, for the first term on the right hand side of (78), we have

$$
\begin{equation*}
G F^{-1 / 2}\left\{\int_{0}^{\infty} e^{-p v} v^{-1 / 2} \mathrm{~d} v-\varepsilon_{1}(p)\right\} \sim\left(p H_{1} / \pi\right)^{-1 / 2} \prod_{j=2}^{k}\left(\left(\tilde{\lambda}_{p 1}-z_{j}\right) /\left(\tilde{\lambda}_{p 1}-\tilde{\lambda}_{p j}\right)\right)^{1 / 2} \tag{79}
\end{equation*}
$$

That is, the ratio of the left to the right hand sides of the above display converges to one as $p, n \rightarrow_{\gamma} \infty$, uniformly in $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, almost surely.

Next, by (77), for an arbitrarily small positive $\tau$, we can choose $\kappa$ so that

$$
\left|q(v)-G F^{-1 / 2} v^{-1 / 2}\right|<\tau|G| F^{-1 / 2} v^{-1 / 2}
$$

for all $v \in(0, \kappa]$ and all $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$. Therefore, $\varepsilon_{2}$ is asymptotically dominated by the first term on the right hand side of (78).

Finally, let

$$
\eta=\inf _{\left[c,--\tilde{\lambda}_{p 2}-\varepsilon\right)}\left\{f(t)-f\left(-\tilde{\lambda}_{p 1}\right)\right\}
$$

Since $f(t)$ is strictly increasing on $\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right], \eta$ is larger than some positive number for any $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, for all sufficiently large $p$, a.s. Therefore,

$$
\left|e^{p f\left(-\tilde{\lambda}_{p 1}\right)} \int_{c}^{-\tilde{\lambda}_{p 2}-\varepsilon} e^{-p f(t)} g(t) \mathrm{d} t\right| \leq e^{-p \theta} \int_{c}^{-\tilde{\lambda}_{p 2}-\varepsilon}|g(t)| \mathrm{d} t,
$$

which is dominated by the right hand side of (79).
Summing up, we have established the following lemma.
Lemma 17. As $p, n \rightarrow \gamma \infty$,

$$
\overline{\mathcal{I}}_{\sigma \varnothing, i n}^{\mathrm{SigD}}=G_{-1}^{\mathrm{SigD}} e^{-p f_{1}\left(\tilde{\lambda}_{p 1}\right)}\left(p H_{1} / \pi\right)^{-1 / 2} \prod_{j=2}^{k}\left(\left(\tilde{\lambda}_{p 1}-z_{j}\right) /\left(\tilde{\lambda}_{p 1}-\tilde{\lambda}_{p j}\right)\right)^{1 / 2}(1+o(1))
$$

where $o(1) \rightarrow 0$ uniformly over $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$ and over $z_{2}, \ldots, z_{k}$ that belong to the (trimmed) domain of integration in (DJO49), a.s.

Repeating the above analysis for the second, third, etc. to the inner-most integral in (DJO49) and combining the results, we obtain Lemma DJO18.
5.6. Proof of Lemma DJO19 about the Laplace approximation for $\overline{\mathcal{I}}_{\sigma_{\varnothing}}^{\mathrm{REG}}$. The proof is very similar to that of Lemma DJO18. As shown in DJO, the inner-most integral in the multiple integral representation (DJO52) of $\overline{\mathcal{I}}_{\sigma_{\varnothing}}^{\mathrm{REG}}$ equals

$$
\begin{equation*}
\overline{\mathcal{I}}_{\sigma_{\varnothing}, i n}^{\mathrm{REG}}=G_{-1}^{\mathrm{REG}} \int_{\tilde{\lambda}_{p 2}+\varepsilon}^{\tilde{\lambda}_{p 1}} e^{-n_{H} f_{2}\left(z_{1}\right)} g_{2}\left(z_{1}\right)(1+o(1)) \mathrm{d} z_{1} \tag{80}
\end{equation*}
$$

First, we will apply the Laplace method to the integral

$$
I(p) \equiv \int_{\tilde{\lambda}_{p 2}+\varepsilon}^{\tilde{\lambda}_{p 1}} e^{-n_{H} f_{2}\left(z_{1}\right)} g_{2}\left(z_{1}\right) \mathrm{d} z_{1}
$$

To line up our analysis with that on pages 81-82 of Olver (1997) and to simplify notations, let us rewrite $I(p)$ as

$$
I(p)=\int_{-\tilde{\lambda}_{p 1}}^{-\tilde{\lambda}_{p 2}-\varepsilon} e^{-n_{H} f(t)} g(t) \mathrm{d} t
$$

where $t=-z_{1}$,

$$
\begin{aligned}
& f(t)=-z_{1+}-a \ln \left(z_{1+}+a\right)+b \ln \left(z_{1+}+b\right)+\frac{1}{2 n_{H}} \sum_{j=k+1}^{p} \ln \left(-t-\tilde{\lambda}_{p j}\right), \text { and } \\
& g(t)=\left(\frac{z_{1+}\left(z_{1+}+a\right)}{z_{1+}^{2}+a \zeta_{1}}\right)^{\frac{1}{2}}\left(\tilde{\lambda}_{p 1}+t\right)^{-1 / 2} \prod_{j=2}^{k}\left(\frac{-t-z_{j}}{-t-\tilde{\lambda}_{p j}}\right)^{1 / 2}
\end{aligned}
$$

Here

$$
\begin{gathered}
z_{j+}=\frac{1}{2}\left\{\zeta_{j}-b+\sqrt{\left(b-\zeta_{j}\right)^{2}+4 a \zeta_{j}}\right\} \text { with } \zeta_{j}=-t s_{j} / 2, \\
a=(N-p) /\left(2 n_{H}\right) \equiv\left(n_{H}+n_{2}-p\right) /\left(2 n_{H}\right), \text { and } b=\left(n_{H}-p\right) /\left(2 n_{H}\right) .
\end{gathered}
$$

Let us show that under the null hypothesis, $f(t)$ has positive continuous derivative for $t$ from the interval $\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right]$ for sufficiently large $p$, a.s. Express $f^{\prime}(t)$ as

$$
f^{\prime}(t)=\varphi^{\prime}(t)+w^{\prime}(t),
$$

where

$$
\varphi(t) \equiv-z_{1+}-a \ln \left(z_{1+}+a\right)+b \ln \left(z_{1+}+b\right) \text { and } w(t) \equiv \frac{1}{2 n_{H}} \sum_{j=k+1}^{p} \ln \left(-t-\tilde{\lambda}_{p j}\right) .
$$

Clearly, function $w(t)$ is decreasing and concave on $t \in\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right]$. We will now show that $\varphi(t)$ is increasing and convex.
Since

$$
z_{1+}^{2}+\left(b-\zeta_{1}\right) z_{1+}-a \zeta_{1}=0,
$$

we have

$$
\begin{equation*}
\zeta_{1}=z_{1+}\left(z_{1+}+b\right) /\left(z_{1+}+a\right) . \tag{81}
\end{equation*}
$$

Therefore, $z_{1+}>\zeta_{1}>0$. Further,

$$
\begin{align*}
\frac{\mathrm{d}}{\mathrm{~d} t} z_{1+} & =\frac{\partial z_{1+}}{\partial \zeta_{1}} \frac{\mathrm{~d} \zeta_{1}}{\mathrm{~d} t}=\frac{1}{2}\left\{1+\frac{\zeta_{1}-b+2 a}{\sqrt{\left(b-\zeta_{1}\right)^{2}+4 a \zeta_{1}}}\right\} \frac{-s_{1}}{2} \\
& =\frac{z_{1+}+a}{2 z_{1+}+b-\zeta_{1}} \frac{-s_{1}}{2}=\frac{\left(a+z_{1+}\right)^{2}}{a b+2 a z_{1+}+z_{1+}^{2}} \frac{-s_{1}}{2}<0 . \tag{82}
\end{align*}
$$

On the other hand,

$$
\begin{equation*}
\frac{\partial}{\partial z_{1+}} \varphi=-\frac{a b+2 a z_{1+}+z_{1+}^{2}}{\left(b+z_{1+}\right)\left(a+z_{1+}\right)}<0 . \tag{83}
\end{equation*}
$$

Therefore, $\varphi^{\prime}(t)$ is positive on $t \in\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right]$. Further,

$$
\frac{\partial^{2}}{\partial z_{1+}^{2}} \varphi=-\frac{b}{\left(z_{1+}+b\right)^{2}}+\frac{a}{\left(z_{1+}+a\right)^{2}}=\frac{(a-b)\left(z_{1+}^{2}-a b\right)}{\left(z_{1+}+b\right)^{2}\left(z_{1+}+a\right)^{2}},
$$

and using (82), we also have

$$
\frac{\mathrm{d}^{2}}{\mathrm{~d} t^{2}} z_{+1}=-\frac{s_{1}^{2}}{2} \frac{a\left(a+z_{1+}\right)^{3}(a-b)}{\left(a b+2 a z_{1+}+z_{1+}^{2}\right)^{3}} .
$$

Hence,

$$
\varphi^{\prime \prime}(t)=\frac{\partial^{2} \varphi}{\partial z_{1+}^{2}}\left(\frac{\mathrm{~d} z_{1+}}{\mathrm{d} t}\right)^{2}+\frac{\partial \varphi}{\partial z_{1+}} \frac{\mathrm{d}^{2} z_{1+}}{\mathrm{d} t^{2}}=\frac{s_{1}^{2}}{4} \frac{\left(a+z_{1+}\right)^{2}(a-b)}{\left(b+z_{1+}\right)^{2}\left(a b+2 a z_{1+}+z_{1+}^{2}\right)}>0 .
$$

The concavity of $w(t)$ and the convexity of $\varphi(t)$ imply that

$$
\begin{equation*}
\min _{t \in\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right]} f^{\prime}(t)>\varphi^{\prime}\left(-\tilde{\lambda}_{p 1}\right)+w^{\prime}\left(-\tilde{\lambda}_{p 2}-\varepsilon\right) . \tag{84}
\end{equation*}
$$

On the other hand, using, first, (82) and (83), and then (81) and the definition of $\zeta_{1}$, we obtain

$$
\varphi^{\prime}(t)=\frac{s_{1}\left(a+z_{1+}\right)}{2\left(b+z_{1+}\right)}=-z_{1+} / t
$$

This and the fact that

$$
\tilde{\lambda}_{p 1} \xrightarrow{\text { a.s. }} \frac{\left(s_{01}+\gamma_{1}\right)\left(s_{01}+1\right) \gamma_{2}}{s_{01}\left(\gamma_{1}+s_{01} \gamma_{2}+\gamma_{2}\right)}
$$

yield, after some algebra,

$$
\begin{equation*}
\varphi^{\prime}\left(-\tilde{\lambda}_{p 1}\right) \xrightarrow{\text { a.s. }} \frac{s_{01}\left(\gamma_{1}+s_{01} \gamma_{2}+\gamma_{2}\right)}{2\left(1+s_{01}\right) \gamma_{2}} . \tag{85}
\end{equation*}
$$

For $w^{\prime}\left(-\tilde{\lambda}_{p 2}-\varepsilon\right)$, we have

$$
w^{\prime}\left(-\tilde{\lambda}_{p 2}-\varepsilon\right) \xrightarrow{\text { a.s. }} \frac{\gamma_{1}}{2} \tilde{m}\left(\lim \tilde{\lambda}_{p 2}+\varepsilon\right) .
$$

Hence, the right hand side of (84) a.s. converges to

$$
\Pi\left(\lim \tilde{\lambda}_{p 2}+\varepsilon, s_{01}\right) \equiv \frac{s_{01}\left(\gamma_{1}+s_{01} \gamma_{2}+\gamma_{2}\right)}{2\left(1+s_{01}\right) \gamma_{2}}+\frac{\gamma_{1}}{2} \tilde{m}\left(\lim \tilde{\lambda}_{p 2}+\varepsilon\right) .
$$

This convergence is uniform in $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$.
Now note that $\Pi\left(y, s_{01}\right)$ is strictly increasing function of $y$ and $s_{01}$ on $y \in\left(\tilde{\beta}_{+}, 1\right]$ and $s_{01}>\bar{s}$. Using the same tools as in the above proof of Lemma DJO18 (more specifically, those used for the analysis of $\Psi\left(y, s_{01}\right)$ ), we find that the limit of $\Pi\left(y, s_{01}\right)$ as $y \rightarrow \tilde{\beta}_{+}$and $s_{01} \rightarrow \bar{s}$ is zero. Therefore, $\Pi\left(\lim \tilde{\lambda}_{p 2}+\varepsilon, s_{01}\right)$ is positive, and thus, by $(84), f^{\prime}(t)$ is positive on $t \in\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right]$. Moreover, there exists a small positive number that is smaller than $f^{\prime}(t)$ for all $t \in\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right]$, all $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, and all sufficiently large $p$, a.s.

Since $f(t)$ is strictly increasing on $\left[-\tilde{\lambda}_{p 1},-\tilde{\lambda}_{p 2}-\varepsilon\right]$, the main contribution to the integral $I(p)$ comes from the vicinity of $-\tilde{\lambda}_{p 1}$. Using (85) and (74), we obtain

$$
f^{\prime}\left(-\tilde{\lambda}_{p 1}\right) \xrightarrow{\text { a.s. }} \gamma_{1} H_{1} .
$$

The rest of the proof is almost identical to that of Lemma DJO18. The only notable difference is that function $g(t)$ has an additional multiplicative term $\left(\frac{z_{1+}\left(z_{1+}+a\right)}{z_{1+}^{2}+a \varsigma_{1}}\right)^{\frac{1}{2}}$. It is straightforward to verify that, at $t=-\tilde{\lambda}_{p 1}$, this term a.s. converges to

$$
\frac{\gamma_{1}+\gamma_{2}+\gamma_{2} s_{01}}{\left(\left(\gamma_{1}+\gamma_{2}+\gamma_{2} s_{01}\right)^{2}-\rho^{2} \gamma_{1}\right)^{1 / 2}},
$$

which explains the presence of the last term in the definition of $\Omega_{j}^{\text {REG }}$ given in Lemma DJO19.
5.7. Proof of Lemma DJO20 about the asymptotic equivalence of $\mathcal{I}_{\sigma_{\varnothing}}$ and $\overline{\mathcal{I}}_{\sigma_{\varnothing}}$. By definition,

$$
\begin{equation*}
\mathcal{I}_{\sigma_{\varnothing}}-\overline{\mathcal{I}}_{\sigma_{\varnothing}}=\sum_{\mathcal{D}} \int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\mathcal{D}_{k-1}} \ldots \int_{\mathcal{D}_{1}}|\omega| \mathcal{F}_{p, n} \prod_{j=1}^{k} \mathrm{~d} z_{j}, \tag{86}
\end{equation*}
$$

where the sum runs over all $\mathcal{D}_{j}$ that are represented by either $\left[\tilde{\lambda}_{j+1}+\varepsilon, \tilde{\lambda}_{j}\right]$ or $\left[\tilde{\lambda}_{j+1}, \tilde{\lambda}_{j+1}+\varepsilon\right]$, and at least one $\mathcal{D}_{j}, j=1, \ldots, k-1$, is represented by $\left[\tilde{\lambda}_{j+1}, \tilde{\lambda}_{j+1}+\varepsilon\right]$. All terms in the above sum can be analyzed symilarly, and here we will focus on the term

$$
J \equiv \int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\tilde{\lambda}_{p k}+\varepsilon}^{\tilde{\lambda}_{p, k-1}} \cdots \int_{\tilde{\lambda}_{p 3}+\varepsilon}^{\tilde{\lambda}_{p 2}} \int_{\tilde{\lambda}_{p 2}}^{\tilde{\lambda}_{p 2}+\varepsilon}|\omega| \mathcal{F}_{p, n} \prod_{j=1}^{k} \mathrm{~d} z_{j} .
$$

As is explained in DJO, we have

$$
\begin{equation*}
|J| \leq \bar{J} \equiv \int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\tilde{\lambda}_{p k}+\varepsilon}^{\tilde{\lambda}_{p, k-1}} \cdots \int_{\tilde{\lambda}_{p 3}+\varepsilon}^{\tilde{\lambda}_{p 2}} \int_{\tilde{\lambda}_{p 2}}^{\tilde{\lambda}_{p 2}+\varepsilon}|\omega| \mathcal{F}_{p, n}\left(\tilde{Z}_{1}\right) \prod_{j=1}^{k} \mathrm{~d} z_{j}, \tag{87}
\end{equation*}
$$

where $\tilde{Z}_{1}=\operatorname{diag}\left\{\tilde{\lambda}_{p 2}+\varepsilon, z_{2}, \ldots, z_{k}\right\}$.
$\underline{\text { Case SigD. Using Lemma DJO16 to approximate } \mathcal{F}_{p, n}\left(\tilde{Z}_{1}\right) \text {, we obtain }}$

$$
\begin{equation*}
\bar{J}=b^{\mathrm{SigD}}(1+o(1)) \int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\tilde{\lambda}_{p k}+\varepsilon}^{\tilde{\lambda}_{p, k-1}} \cdots \int_{\tilde{\lambda}_{p 3}+\varepsilon}^{\tilde{\lambda}_{p 2}} \int_{\tilde{\lambda}_{p 2}}^{\tilde{\lambda}_{p 2}+\varepsilon} e^{-p f_{1}\left(\tilde{z}_{1}\right)} \tilde{g}_{1} \mathrm{~d} z_{1} \prod_{j=2}^{k}\left\{e^{-p f_{j}\left(z_{j}\right)} g_{j} \mathrm{~d} z_{j}\right\} \tag{88}
\end{equation*}
$$

where $o(1) \rightarrow 0$ uniformly in $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, and the definitions of $f_{j}, g_{j}, \tilde{z}_{1}$, and $\tilde{g}_{1}$ are summarized in Table 1 .

Table 1
Definition of $f_{j}, g_{j}, \tilde{z}_{1}$, and $\tilde{g}_{1}$ used in equation (88).

| Quantity | Definition |
| :--- | :--- |
| $f_{j}(z)$ | $=\frac{N-p+2}{2 p} \ln \left(1-\frac{s_{j} z}{1+s_{j}}\right)+\frac{1}{2 p} \sum_{i=k+1}^{p} \ln \left(z-\tilde{\lambda}_{p i}\right), j=1, \ldots, k$, |
| $g_{j}$ | $=\prod_{s=1}^{j}\left(\tilde{\lambda}_{p s}-z_{j}\right)^{-1 / 2} \prod_{i=j+1}^{k}\left(\left(z_{j}-z_{i}\right) /\left(z_{j}-\tilde{\lambda}_{p i}\right)\right)^{1 / 2}, j=2, \ldots, k$, |
| $\tilde{z}_{1}$ | $=\tilde{\lambda}_{p 2}+\varepsilon$, |
| $\tilde{g}_{1}$ | $=\left(\tilde{\lambda}_{p 1}-\tilde{z}_{1}\right)^{-1 / 2} \prod_{i=2}^{k}\left\{\left(\frac{z_{1}-z_{i}}{\tilde{z}_{1}-\lambda_{p i}}\right)^{1 / 2} \frac{z_{1}-z_{i}}{z_{1}-z_{i}}\right\} \prod_{s=1}^{p}\left\|\frac{z_{1}-\lambda_{p s}}{z_{1}-\lambda_{p s}}\right\|^{-1 / 2}$. |

As follows from the proof of Lemma DJO18, $f_{j}(z), j=1, \ldots, k$ are strictly decreasing functions for $z$ from the integration domain of $J$. Therefore, we have

$$
\begin{aligned}
\bar{J} & \leq b^{\mathrm{SigD}}(1+o(1)) \int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\tilde{\lambda}_{p k}+\varepsilon}^{\tilde{\lambda}_{p, k-1}} \cdots \int_{\tilde{\lambda}_{p 3}+\varepsilon}^{\tilde{\lambda}_{p 2}} \int_{\tilde{\lambda}_{p 2}}^{\tilde{\lambda}_{p 2}+\varepsilon} e^{-p f_{1}\left(\tilde{z}_{1}\right)} \tilde{g}_{1} \mathrm{~d} z_{1} \prod_{j=2}^{k}\left\{e^{-p f_{j}\left(\tilde{\lambda}_{p j}\right)} g_{j} \mathrm{~d} z_{j}\right\} \\
& \leq C^{\operatorname{SigD}} e^{-p f_{1}\left(\tilde{\lambda}_{p 2}+\varepsilon\right)} \prod_{j=2}^{k} e^{-p f_{j}\left(\tilde{\lambda}_{p j}\right)}
\end{aligned}
$$

where $C^{\mathrm{SigD}}$ is a positive quantity that depends only on $s_{01}, \ldots, s_{0 k}$.
In the proof of Lemma DJO18, we have seen that not only $f_{1}(z)$ is strictly decreasing in $z$ on $z \in\left[\tilde{\lambda}_{p 2}, \tilde{\lambda}_{p 1}\right]$, but also there exists a fixed negative number such that the derivative of $f_{1}(z)$ is smaller than that number. Therefore, there exists $\eta>0$, such that $f_{1}\left(\tilde{\lambda}_{p 2}+\varepsilon\right)>f_{1}\left(\tilde{\lambda}_{p 1}\right)+\eta$ and

$$
\bar{J} \leq C^{\mathrm{SigD}} e^{-p \eta} \prod_{j=1}^{k} e^{-p f_{j}\left(\tilde{\lambda}_{p j}\right)}
$$

On the other hand, by Lemma DJO18, the right hand side of the displayed inequality is asymptotically dominated by $\overline{\mathcal{I}}_{\sigma_{\varnothing}}$. Repeating the above arguments for the other components of $\mathcal{I}_{\sigma_{\varnothing}}-\overline{\mathcal{I}}_{\sigma_{\varnothing}}$ (that is, the components of the sum in (86) other than $J$ ), we establish Lemma DJO20 for the SigD case.

Case REG. A proof of this case is similar to that for SigD. Using Lemma DJO17 to approximate $\mathcal{F}_{p, n}\left(\tilde{Z}_{1}\right)$, we obtain

$$
\begin{equation*}
\bar{J}=b^{\mathrm{REG}}(1+o(1)) \int_{\tilde{x}_{0}}^{\tilde{\lambda}_{p k}} \int_{\tilde{\lambda}_{p k}+\varepsilon}^{\tilde{\lambda}_{p, k-1}} \cdots \int_{\tilde{\lambda}_{p 3}+\varepsilon}^{\tilde{\lambda}_{p 2}} \int_{\tilde{\lambda}_{p 2}}^{\tilde{\lambda}_{p 2}+\varepsilon} e^{-n_{H} f_{1}\left(\tilde{z}_{1}\right)} \tilde{g}_{1} \mathrm{~d} z_{1} \prod_{j=2}^{k}\left\{e^{-n_{H} f_{j}\left(z_{j}\right)} g_{j} \mathrm{~d} z_{j}\right\} \tag{89}
\end{equation*}
$$

where $o(1) \rightarrow 0$ uniformly in $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, and the definitions of $f_{j}, g_{j}, \tilde{z}_{1}$, and $\tilde{g}_{1}$ are summarized in Table 2.

Proceeding exactly as in SigD case, we obtain inequality

$$
\bar{J} \leq C^{\mathrm{REG}} e^{-n_{H} \eta} \prod_{j=1}^{k} e^{-n_{H} f_{j}\left(\tilde{\lambda}_{p j}\right)}
$$

so Lemma DJO20 follows by a similar argument.
Table 2
Definition of $f_{j}, g_{j}, \tilde{z}_{1}$, and $\tilde{g}_{1}$ used in equation (89).

| Quantity | Definition |
| :--- | :--- |
| $f_{j}\left(z_{j}\right)$ | $=-z_{j+}-a \ln \left(z_{j+}+a\right)+b \ln \left(z_{j+}+b\right)+\frac{1}{2 n_{H}} \sum_{i=k+1}^{p} \ln \left(z_{j}-\tilde{\lambda}_{p i}\right), j=2, \ldots, k$, |
| $g_{j}$ | $=\left(\frac{z_{j+}\left(z_{j+}+a\right)}{z_{j+}^{2}+a \zeta_{j}}\right)^{1 / 2} \prod_{s=1}^{j}\left(\tilde{\lambda}_{p s}-z_{j}\right)^{-1 / 2} \prod_{i=j+1}^{k}\left(\frac{z_{j}-z_{i}}{z_{j}-\lambda_{p i}}\right)^{1 / 2}, j=2, \ldots, k$, |
| $\tilde{z}_{1}$ | $=\tilde{\lambda}_{p 2}+\varepsilon$, |
| $f_{1}\left(\tilde{z}_{1}\right)$ | $=-\tilde{z}_{1+}-a \ln \left(\tilde{z}_{1+}+a\right)+b \ln \left(\tilde{z}_{1+}+b\right)+\frac{1}{2 n_{H}} \sum_{i=k+1}^{p} \ln \left(\tilde{z}_{1}-\tilde{\lambda}_{p i}\right)$, |
| $\tilde{g}_{1}$ | $=\left(\frac{\tilde{z}_{1+}\left(\tilde{z}_{1+}+a\right)}{\tilde{z}_{1+}^{2}+a \tilde{\zeta}_{1}}\right)^{1 / 2}\left(\tilde{\lambda}_{p 1}-\tilde{z}_{1}\right)^{-1 / 2} \prod_{i=2}^{k}\left\{\left(\frac{\tilde{z}_{1}-z_{i}}{\tilde{z}_{1}-\lambda_{p i}}\right)^{1 / 2} \frac{z_{1}-z_{i}}{\tilde{z}_{1}-z_{i}}\right\} \prod_{s=1}^{p}\left\|\frac{z_{1}-\tilde{\lambda}_{p s}}{\tilde{z}_{1}-\lambda_{p s}}\right\|^{-1 / 2}$. |
| $\tilde{z}_{1+}$ | $=\frac{1}{2}\left\{\tilde{\zeta}_{1}-b+\sqrt{\left(b-\tilde{\zeta}_{1}\right)^{2}+4 a \tilde{\zeta}_{1}}\right\}$, |
| $\tilde{\zeta}_{1}$ | $=\tilde{z}_{1} s_{1} / 2$. |

5.8. Proof of Lemma DJO21 about the asymptotic negligibility of $\mathcal{I}_{\sigma_{\tau}}$ with $\tau \neq \varnothing$. We consider here only the case of $\tau=1,2, \ldots, k$. The analysis for the other subsets $\tau \subseteq 1,2, \ldots, k$ is very similar.

As shown in DJO,

$$
\left|\mathcal{I}_{\sigma_{\tau}}\right| \leq \frac{1}{2^{k}} \int_{\mathcal{K}_{k+1}} \ldots \int_{\mathcal{K}_{k+1}}|\omega| \mathcal{F}_{p, n}\left(Z_{\eta}\right) \prod_{j=1}^{k}\left|\mathrm{~d} z_{j}\right|
$$

where $Z_{\eta}=\operatorname{diag}\left\{\tilde{x}_{0}+k \eta, \ldots, \tilde{x}_{0}+2 \eta, \tilde{x}_{0}+\eta\right\}$ and $\eta$ is a fixed small positive number. Function $\mathcal{F}_{p, n}\left(Z_{\eta}\right)$ can now be approximated using Lemmas DJO16 and DJO17.

Case SigD. Using Lemma DJO16, we obtain

$$
\begin{equation*}
\left|\mathcal{I}_{\sigma_{\tau}}\right| \leq C \prod_{j=1}^{k}\left(1-\frac{s_{j}\left(\tilde{x}_{0}+j \eta\right)}{1+s_{j}}\right)^{-\frac{N-p+2}{2}} \int_{\mathcal{K}_{k+1}} \ldots \int_{\mathcal{K}_{k+1}}|\omega| \prod_{j=1}^{k}\left|\mathrm{~d} z_{j}\right|, \tag{90}
\end{equation*}
$$

where $C$ is some positive constant. The above inequality holds uniformly in $s=s_{0}+\delta / \sqrt{p}$ with $\left(\delta_{1}, \ldots, \delta_{k}\right)$ from a compact subset of $\mathbb{R}^{k}$, for sufficiently large $p$, a.s.

Using the definitions of $\omega$ and $\mathcal{K}_{k+1}$, we obtain

$$
\begin{equation*}
\int_{\mathcal{K}_{k+1}} \ldots \int_{\mathcal{K}_{k+1}}|\omega| \prod_{j=1}^{k}\left|\mathrm{~d} z_{j}\right| \leq C_{1} \prod_{j=1}^{k} \prod_{s=k+1}^{p}\left(\tilde{x}_{0}-\tilde{\lambda}_{p s}\right)^{-1 / 2} \tag{91}
\end{equation*}
$$

where $C_{1}$ is a positive constant. Combining this inequality with (90), we obtain

$$
\left|\mathcal{I}_{\sigma_{\tau}}\right| \leq C_{2} \prod_{j=1}^{k}\left[\left(1-\frac{s_{j}\left(\tilde{x}_{0}+j \eta\right)}{1+s_{j}}\right)^{-\frac{N-p+2}{2}} \prod_{s=k+1}^{p}\left(\tilde{x}_{0}-\tilde{\lambda}_{p s}\right)^{-1 / 2}\right] .
$$

Now recall (from the proof of Lemma DJO18) that functions

$$
f_{j}(z)=\frac{N-p+2}{2 p} \ln \left(1-\frac{s_{j} z}{1+s_{j}}\right)+\frac{1}{2 p} \sum_{j=k+1}^{p} \ln \left(z-\tilde{\lambda}_{p s}\right), j=1, \ldots, k,
$$

are decreasing on $z \in\left[\tilde{x}_{0}, 1\right]$ and have there derivatives that are bounded away from zero, for sufficiently large $p$, a.s. This implies that there exists a small positive $\eta_{1}$ such that

$$
\left|\mathcal{I}_{\sigma_{\tau}}\right| \leq C_{3} e^{-p \eta_{1}} \prod_{j=1}^{k} e^{-p f_{j}\left(\tilde{\lambda}_{p j}\right)},
$$

and therefore, by Lemma DJO18, $\mathcal{I}_{\sigma_{\tau}}$ is asymptotically dominated by $\mathcal{I}_{\sigma_{\varnothing}}$.
Case REG. Using Lemma DJO17, we obtain

$$
\begin{equation*}
\left|\mathcal{I}_{\sigma_{\tau}}\right| \leq C \prod_{j=1}^{k} e^{-n_{H} \varphi_{j}\left(\tilde{x}_{0}+j \eta\right)} \int_{\mathcal{K}_{k+1}} \ldots \int_{\mathcal{K}_{k+1}}|\omega| \prod_{j=1}^{k}\left|\mathrm{~d} z_{j}\right|, \tag{92}
\end{equation*}
$$

where $C$ is some positive constant and

$$
\begin{aligned}
& \varphi_{j}(z)=-z_{+}-a \ln \left(z_{+}+a\right)+b \ln \left(z_{+}+b\right) \\
& z_{+}=\frac{1}{2}\left\{\zeta-b+\sqrt{(b-\zeta)^{2}+4 a \zeta}\right\} \text { with } \zeta=z s_{j} / 2 \\
& \text { imsart-aos ver. 2012/04/10 file: supercriticalsm1.tex date: January } 28,2017
\end{aligned}
$$

Using (91) in (92), we obtain

$$
\left|\mathcal{I}_{\sigma_{\tau}}\right| \leq C_{2} \prod_{j=1}^{k}\left[e^{-n_{H} \varphi_{j}\left(\tilde{x}_{0}+j \eta\right)} \prod_{s=k+1}^{p}\left(\tilde{x}_{0}-\tilde{\lambda}_{p s}\right)^{-1 / 2}\right]
$$

Now recall (from the proof of Lemma DJO19) that functions

$$
f_{j}(z)=\varphi_{j}(z)+\frac{1}{2 n_{H}} \sum_{j=k+1}^{p} \ln \left(z-\tilde{\lambda}_{p s}\right), j=1, \ldots, k
$$

are decreasing on $z \in\left[\tilde{x}_{0}, 1\right]$ and have there derivatives that are bounded away from zero, for sufficiently large $p$, a.s. This implies that there exists a small positive $\eta_{1}$ such that

$$
\left|\mathcal{I}_{\sigma_{\tau}}\right| \leq C_{3} e^{-p \eta_{1}} \prod_{j=1}^{k} e^{-p f_{j}\left(\tilde{\lambda}_{p j}\right)}
$$

and therefore, by Lemma $\operatorname{DJO} 19, \mathcal{I}_{\sigma_{\tau}}$ is asymptotically dominated by $\mathcal{I}_{\sigma_{\varnothing}}$.

## 6. Local Asymptotic Normality.

6.1. Proof of Theorem DJO23 about the quadratic approximation of the log likelihood ratio. Case SigD. Theorem DJO22 yields

$$
\begin{equation*}
L_{p}^{\mathrm{SigD}}(\theta, \Lambda)=(1+o(1)) \frac{k_{p, n}^{\mathrm{SigD}}(s)}{k_{p, n}^{\mathrm{SigD}}\left(s_{0}\right)} \prod_{j=1}^{k}\left(\frac{1-\frac{s_{j}}{1+s_{j}} \tilde{\lambda}_{p j}}{1-\frac{s_{j}}{1+s_{0 j}} \tilde{\lambda}_{p j}}\right)^{\frac{p-N-2}{2}} \tag{93}
\end{equation*}
$$

The right hand side does not depend on $b^{\operatorname{SigD}}(s) / b^{\operatorname{SigD}}\left(s_{0}\right)$ because the latter ratio is asymptotically equivalent to one. Taking logarithm of both sides of (93) and simplifying, we obtain

$$
\begin{equation*}
\ln L_{p}^{\mathrm{SigD}}(\theta, \Lambda)=\sum_{j=1}^{k}\left\{\frac{n_{2}}{2} \ln \frac{1+s_{j}}{1+s_{0 j}}-\frac{p}{2} \ln \frac{s_{j}}{s_{0 j}}-\frac{N-p}{2} \ln \frac{1+s_{j}-s_{j} \tilde{\lambda}_{p j}}{1+s_{0 j}-s_{0 j} \tilde{\lambda}_{p j}}\right\}+o(1) \tag{94}
\end{equation*}
$$

Using the identity

$$
\begin{equation*}
\tilde{\lambda}_{p j}=\frac{\alpha_{n} \lambda_{p j}}{1+\alpha_{n} \lambda_{p j}} \text { with } \alpha_{n}=n_{H} / n_{2} \tag{95}
\end{equation*}
$$

we rewrite (94) as

$$
\begin{equation*}
\ln L_{p}^{\mathrm{SigD}}(\theta, \Lambda)=\sum_{j=1}^{k}\left\{\frac{n_{2}}{2} \ln \frac{1+s_{j}}{1+s_{0 j}}-\frac{p}{2} \ln \frac{s_{j}}{s_{0 j}}-\frac{N-p}{2} \ln \frac{1+s_{j}+\alpha_{n} \lambda_{p j}}{1+s_{0 j}+\alpha_{n} \lambda_{p j}}\right\}+o(1) \tag{96}
\end{equation*}
$$

Expanding the logarithms in (96) in the powers of $\theta p^{-1 / 2}$ up to and including terms $\theta^{2} p^{-1}$, we obtain

$$
\begin{aligned}
\ln L_{p}^{\mathrm{SigD}}(\theta, \Lambda)=\sum_{j=1}^{k} \quad & \left\{\frac{n_{2}}{2}\left(\theta_{j} p^{-1 / 2} \frac{w\left(s_{0 j}\right)}{1+s_{0 j}}-\frac{1}{2} \theta_{j}^{2} p^{-1} \frac{w^{2}\left(s_{0 j}\right)}{\left(1+s_{0 j}\right)^{2}}+o\left(p^{-1}\right)\right)\right. \\
& -\frac{p}{2}\left(\theta_{j} p^{-1 / 2} \frac{w\left(s_{0 j}\right)}{s_{0 j}}-\frac{1}{2} \theta_{j}^{2} p^{-1} \frac{w^{2}\left(s_{0 j}\right)}{s_{0 j}^{2}}+o\left(p^{-1}\right)\right) \\
& \left.-\frac{N-p}{2}\left(\theta_{j} p^{-1 / 2} \frac{w\left(s_{0 j}\right)}{1+s_{0 j}+\alpha_{n} \lambda_{p j}}-\frac{1}{2} \theta_{j}^{2} p^{-1} \frac{w^{2}\left(s_{0 j}\right)}{\left(1+s_{0 j}+\alpha_{n} \lambda_{p j}\right)^{2}}+o\left(p^{-1}\right)\right)\right\} \\
& +o(1)
\end{aligned}
$$

Consider, first, the terms linear in $\theta_{j}$. They can be rewritten as

$$
T_{\theta_{j}} \equiv \frac{1}{2} \theta_{j} p^{1 / 2} w\left(s_{0 j}\right)\left(\frac{1}{c_{2}} \frac{1}{1+s_{0 j}}-\frac{1}{s_{0 j}}-\frac{r^{2}}{c_{1} c_{2}} \frac{1}{1+s_{0 j}+\alpha_{n} \lambda_{p j}}\right) .
$$

Expanding the last term in the brackets around $\lambda_{p j}=l_{j}$ and using the fact that, by Theorem DJO10, $\sqrt{p}\left(\lambda_{p j}-l_{j}\right)=O_{\mathrm{P}}(1)$, we get

$$
T_{\theta_{j}}=\frac{1}{2} \theta_{j} p^{1 / 2} w\left(s_{0 j}\right)\left(\frac{1}{c_{2}} \frac{1}{1+s_{0 j}}-\frac{1}{s_{0 j}}-\frac{r^{2}}{c_{1} c_{2}} \frac{1}{1+s_{0 j}+\alpha_{n} l_{j}}+\frac{r^{2}}{c_{1}^{2}} \frac{\lambda_{p j}-l_{j}}{\left(1+s_{0 j}+\alpha_{n} l_{j}\right)^{2}}\right)+o_{\mathrm{P}}(1) .
$$

Simplifying this using identities

$$
\begin{align*}
\alpha_{n} & =c_{2} / c_{1}, \\
r^{2} & =c_{1}+c_{2}-c_{1} c_{2}, \text { and }  \tag{97}\\
l_{j} & =\left(s_{0 j}+c_{1}\right)\left(s_{0 j}+1\right) /\left(\left(1-c_{2}\right) s_{0 j}-c_{2}\right),
\end{align*}
$$

we obtain

$$
T_{\theta_{j}}=\theta_{j} \sqrt{p}\left(\lambda_{p j}-l_{j}\right) w\left(s_{0 j}\right) \frac{\left(\left(1-c_{2}\right) s_{0 j}-c_{2}\right)^{2}}{2 r^{2}\left(1+s_{0 j}\right)^{2} s_{0 j}^{2}}+o_{\mathrm{P}}(1) .
$$

On the other hand, for SigD,

$$
\begin{equation*}
w\left(s_{0 j}\right)=2\left(\rho s_{0 j}\left(1+s_{0 j}\right) /\left(\left(1-\gamma_{2}\right) s_{0 j}-\gamma_{2}\right)\right)^{2}, \tag{98}
\end{equation*}
$$

which yields

$$
\begin{equation*}
T_{\theta_{j}}=\theta_{j} \sqrt{p}\left(\lambda_{p j}-l_{j}\right)+o_{\mathrm{P}}(1) . \tag{99}
\end{equation*}
$$

Now, for the quadratic terms in $\theta_{j}$, we have

$$
T_{\theta_{j}^{2}} \equiv \frac{1}{4} \theta_{j}^{2} w^{2}\left(s_{0 j}\right)\left(-\frac{1}{c_{2}\left(1+s_{0 j}\right)^{2}}+\frac{1}{s_{0 j}^{2}}+\frac{r^{2}}{c_{1} c_{2}\left(1+s_{0 j}+\alpha_{n} \lambda_{p j}\right)^{2}}\right) .
$$

Using identities (97) and the fact that $\lambda_{p j}-l_{j}=o_{\mathrm{P}}(1)$, we obtain

$$
\begin{equation*}
T_{\theta_{j}^{2}}=-\frac{1}{4} \theta_{j}^{2} w^{2}\left(s_{0 j}\right) \frac{\left(1-c_{2}\right) s_{0 j}^{2}-2 c_{2} s_{0 j}-c_{1}-c_{2}}{\left(1+s_{0 j}\right)^{2} s_{0 j}^{2} r^{2}}+o_{\mathrm{P}}(1) . \tag{100}
\end{equation*}
$$

Recall that the asymptotic variance of $\lambda_{p j}$ equals $\tau^{2}\left(s_{0 j}\right)=w\left(s_{0 j}\right) t\left(s_{0 j}\right)$, (see Theorem DJO10), where

$$
t\left(s_{0 j}\right) \equiv\left(1-\gamma_{2}\right)\left(s_{0 j}-\underline{\mathrm{s}}\right)\left(s_{0 j}-\bar{s}\right) / d^{2}\left(s_{0 j}\right)=\frac{\left(1-\gamma_{2}\right) s_{0 j}^{2}+2 \gamma_{2} s_{0 j}-\gamma_{1}-\gamma_{2}}{\left(\left(1-\gamma_{2}\right) s_{0 j}-\gamma_{2}\right)^{2}},
$$

Using this with (98) in (100), we obtain

$$
\begin{equation*}
T_{\theta_{j}^{2}}=-\frac{1}{2} \theta_{j}^{2} \tau^{2}\left(s_{0 j}\right)+o_{\mathrm{P}}(1) \tag{101}
\end{equation*}
$$

Combining (99) and (101) yields

$$
\ln L_{p}^{\mathrm{SigD}}(\theta, \Lambda)=\sum_{j=1}^{k}\left\{\theta_{j} \sqrt{p}\left(\lambda_{p j}-l_{j}\right)-\frac{1}{2} \theta_{j}^{2} \tau^{2}\left(s_{0 j}\right)\right\}+o_{\mathrm{P}}(1) .
$$

It remains to note that, by construction, the above $o_{\mathrm{P}}(1)$ term is uniform in $\left(\theta_{1}, \ldots, \theta_{k}\right)$ from any compact subset of $\mathbb{R}^{k}$.

Case REG. Theorem DJO22 yields

$$
\begin{equation*}
L_{p}^{\mathrm{REG}}(\theta, \Lambda)=(1+o(1)) \frac{k_{p, n}^{\mathrm{REG}}(s)}{k_{p, n}^{\mathrm{REG}}\left(s_{0}\right)} \prod_{j=1}^{k} e^{-n_{H}\left(\varphi\left(\bar{z}_{j+}\right)-\varphi\left(\bar{z}_{0 j+}\right)\right)}, \tag{102}
\end{equation*}
$$

where

$$
\begin{aligned}
\bar{z}_{j+} & =\frac{1}{2}\left\{\tilde{\lambda}_{p j} s_{j} / 2-b+\sqrt{\left(b-\tilde{\lambda}_{p j} s_{j} / 2\right)^{2}+2 a \tilde{\lambda}_{p j} s_{j}}\right\} \\
\bar{z}_{0 j+} & =\frac{1}{2}\left\{\tilde{\lambda}_{p j} s_{0 j} / 2-b+\sqrt{\left(b-\tilde{\lambda}_{p j} s_{0 j} / 2\right)^{2}+2 a \tilde{\lambda}_{p j} s_{0 j}}\right\}
\end{aligned}
$$

and

$$
\varphi(z)=-z-a \ln (z+a)+b \ln (z+b) .
$$

Taking logarithm of both sides of (102) and simplifying, we obtain

$$
\begin{equation*}
\ln L_{p}^{\mathrm{REG}}(\theta, \Lambda)=\sum_{j=1}^{k}\left\{-\frac{n_{H}}{2}\left(s_{j}-s_{0 j}\right)-\frac{p}{2} \ln \frac{s_{j}}{s_{0 j}}-n_{H}\left(\varphi\left(\bar{z}_{j+}\right)-\varphi\left(\bar{z}_{0 j+}\right)\right)\right\}+o(1) \tag{103}
\end{equation*}
$$

Expanding the difference $\bar{z}_{j+}-\bar{z}_{0 j+}$ in the powers of $\theta p^{-1 / 2}$ up to and including terms $\theta^{2} p^{-1}$, we obtain

$$
\bar{z}_{j+}-\bar{z}_{0 j+}=\frac{1}{2} R_{0}^{-1} \tilde{\lambda}_{p j}\left(\bar{z}_{0 j+}+a\right) \theta_{j} p^{-1 / 2} w\left(s_{0 j}\right)+\frac{a(b-a)}{4} R_{0}^{-3} \tilde{\lambda}_{p j}^{2} \theta_{j}^{2} p^{-1} w^{2}\left(s_{0 j}\right)+o_{\mathrm{P}}\left(p^{-1}\right)
$$

where

$$
R_{0}=\sqrt{\left(b-\tilde{\lambda}_{p j} s_{0 j} / 2\right)^{2}+2 a \tilde{\lambda}_{p j} s_{0 j}} .
$$

On the other hand,

$$
\ln \frac{\bar{z}_{j+}+a}{\bar{z}_{0 j+}+a}=\frac{\bar{z}_{j+}-\bar{z}_{0 j+}}{\bar{z}_{0 j+}+a}-\frac{\left(\bar{z}_{j+}-\bar{z}_{0 j+}\right)^{2}}{2\left(\bar{z}_{0 j+}+a\right)^{2}}+o_{\mathrm{P}}\left(p^{-1}\right)
$$

and

$$
\ln \frac{\bar{z}_{j+}+b}{\bar{z}_{0 j+}+b}=\frac{\bar{z}_{j+}-\bar{z}_{0 j+}}{\bar{z}_{0 j+}+b}-\frac{\left(\bar{z}_{j+}-\bar{z}_{0 j+}\right)^{2}}{2\left({\left(\bar{z}_{0 j+}+b\right)^{2}}^{2}\right.}+o_{\mathrm{P}}\left(p^{-1}\right)
$$

Therefore,

$$
\begin{align*}
\varphi\left(\bar{z}_{0 j+}\right)-\varphi\left(\bar{z}_{j+}\right)= & \frac{1}{2} R_{0}^{-1} \tilde{\lambda}_{p j} \theta_{j} p^{-1 / 2} w\left(s_{0 j}\right) \frac{\bar{z}_{0 j+}^{2}+2 a \bar{z}_{0 j+}+a b}{\bar{z}_{0 j+}+b} \\
& +\frac{a(b-a)}{4} R_{0}^{-3} \tilde{\lambda}_{p j}^{2} \theta_{j}^{2} p^{-1} w^{2}\left(s_{0 j}\right) \frac{\bar{z}_{0 j+}^{2}+2 a \bar{z}_{0 j+}+a b}{\left(\bar{z}_{0 j+}+a\right)\left(\bar{z}_{0 j+}+b\right)}  \tag{104}\\
& +\frac{b-a}{8} R_{0}^{-2} \tilde{\lambda}_{p j}^{2} \theta_{j}^{2} p^{-1} w^{2}\left(s_{0 j}\right) \frac{\bar{z}_{0 j+}^{2}-a b}{\left(\bar{z}_{0 j+}+b\right)^{2}}+o_{\mathrm{P}}\left(p^{-1}\right)
\end{align*}
$$

By Theorem DJO4, as $p, n \rightarrow \gamma$,

$$
\tilde{\lambda}_{p j} \xrightarrow{a . s .} \frac{\gamma_{2}\left(s_{0 j}+\gamma_{1}\right)\left(s_{0 j}+1\right)}{s_{0 j}\left(\left(\gamma_{1}+\gamma_{2}\left(1+s_{0 j}\right)\right)\right.},
$$

Further, by definition, $b \rightarrow \frac{1}{2}\left(1-\gamma_{1}\right)$ and $a \rightarrow \frac{1}{2}\left(1+\gamma_{1} / \gamma_{2}-\gamma_{1}\right)$. These convergences imply that

$$
\begin{aligned}
R_{0} & \xrightarrow{\text { a.s. }} \frac{1}{2} \frac{\gamma_{2} s_{0 j}^{2}+2 s_{0 j}\left(\gamma_{1}+\gamma_{2}\right)+\gamma_{1}^{2}+\gamma_{1}+\gamma_{2}}{\gamma_{1}+\gamma_{2}\left(1+s_{0 j}\right)}, \text { and } \\
\bar{z}_{0 j+} & \xrightarrow{\text { a.s. }} \frac{1}{2}\left(s_{0 j}+\gamma_{1}\right) .
\end{aligned}
$$

Using these results, we can simplify the last two lines of (104) to obtain

$$
\begin{align*}
\varphi\left(\bar{z}_{0 j+}\right)-\varphi\left(\bar{z}_{j+}\right)= & \frac{1}{2} R_{0}^{-1} \tilde{\lambda}_{p j} \theta_{j} p^{-1 / 2} w\left(s_{0 j}\right) \frac{\bar{z}_{j+}^{2}+2 a \bar{z}_{0 j+}+a b}{\bar{z}_{0 j+}+b} \\
& -\frac{\gamma_{1}\left(s_{0 j}+\gamma_{1}\right)^{2} \theta_{j}^{2} p^{-1} w^{2}\left(s_{0 j}\right)}{4 s_{0 j}^{2}\left(\gamma_{2} s_{0 j}^{2}+2 s_{0 j}\left(\gamma_{1}+\gamma_{2}\right)+\gamma_{1}^{2}+\gamma_{1}+\gamma_{2}\right)}  \tag{105}\\
& +o_{\mathrm{P}}\left(p^{-1}\right) .
\end{align*}
$$

For the first two terms in the figure brackets in (103), we have

$$
\begin{equation*}
-\frac{1}{2}\left(s_{j}-s_{0 j}\right)-\frac{c_{1}}{2} \ln \frac{s_{j}}{s_{0 j}}=-\frac{s_{0 j}+c_{1}}{2 s_{0 j}} \theta_{j} p^{-1 / 2} w\left(s_{0 j}\right)+\frac{c_{1}}{4} \frac{\theta_{j}^{2} p^{-1} w^{2}\left(s_{0 j}\right)}{s_{0 j}^{2}}+o\left(p^{-1}\right) . \tag{106}
\end{equation*}
$$

Using (105) and (106) in (103) yields

$$
\begin{align*}
\ln L_{p}^{\mathrm{REG}}(\theta, \Lambda)=\sum_{j=1}^{k} \quad & \left\{-n_{H} \frac{s_{0 j}+c_{1}}{2 s_{0 j}} \theta_{j} p^{-1 / 2} w\left(s_{0 j}\right)\right. \\
& +n_{H} \frac{1}{2} R_{0}^{-1} \tilde{\lambda}_{p j} \theta_{j} p^{-1 / 2} w\left(s_{0 j}\right) \frac{\bar{z}_{0 j+}^{2}+2 a \bar{z}_{0 j+}+a b}{\bar{z}_{0 j+}+b}  \tag{107}\\
& \left.-\frac{1}{2} \theta_{j}^{2} \tau^{2}\left(s_{0 j}\right)\right\}+o_{\mathrm{P}}(1) .
\end{align*}
$$

Finally, expand the second line of (107) in the powers of $\lambda_{p j}-l_{j}$ up to and including linear terms. To derive such an expansion, note that

$$
\begin{aligned}
\tilde{\lambda}_{p j} & =\frac{c_{2}\left(s_{0 j}+c_{1}\right)\left(s_{0 j}+1\right)}{s_{0 j}\left(c_{1}+c_{2}+s_{0 j} c_{2}\right)}+\frac{c_{1} c_{2}\left(-s_{0 j}+c_{2}+s_{0 j} c_{2}\right)^{2}}{s_{0 j}^{2}\left(c_{1}+c_{2}+s_{0 j} c_{2}\right)^{2}}\left(\lambda_{p j}-l_{j}\right)+o\left(\lambda_{p j}-l_{j}\right) \\
\bar{z}_{0 j+} & =\frac{s_{0 j}+c_{1}}{2}+\frac{c_{1}\left(c_{2}+s_{0 j} c_{2}-s_{0 j}\right)^{2}}{2 s_{0 j}\left(c_{1}+c_{2}+c_{1}^{2}+2 s_{0 j}\left(c_{1}+c_{2}\right)+s_{0 j}^{2} c_{2}\right)}\left(\lambda_{p j}-l_{j}\right)+o\left(\lambda_{p j}-l_{j}\right),
\end{aligned}
$$

and

$$
R_{0}=2 \bar{z}_{0 j+}+\frac{1}{2}\left(1-c_{1}\right)-\frac{s_{0 j}}{2} \tilde{\lambda}_{p j} .
$$

Using these equations, we obtain, after some algebra,

$$
\begin{aligned}
R_{0}^{-1} \tilde{\lambda}_{p j} \frac{\bar{z}_{0 j+}^{2}+2 a \bar{z}_{0 j+}+a b}{\bar{z}_{0 j+}+b}= & \frac{s_{0 j}+c_{1}}{s_{0 j}}+\frac{c_{1}\left(c_{2}+s_{0 j} c_{2}-s_{0 j}\right)^{2}}{s_{0 j}^{2}\left(c_{1}+c_{2}+c_{1}^{2}+2 s_{0 j}\left(c_{1}+c_{2}\right)+s_{0 j}^{2} c_{2}\right)}\left(\lambda_{p j}-l_{j}\right) \\
& +o\left(\lambda_{p j}-l_{j}\right) \\
= & \frac{s_{0 j}+c_{1}}{s_{0 j}}+\frac{2 c_{1}}{w\left(s_{0 j}\right)}\left(\lambda_{p j}-l_{j}\right)+o_{\mathrm{P}}\left(p^{-1 / 2}\right) .
\end{aligned}
$$

Using this equality in the second line of (107) and simplifying, we obtain

$$
\ln L_{p}^{\mathrm{REG}}(\theta, \Lambda)=\sum_{j=1}^{k}\left\{\theta_{j} \sqrt{p}\left(\lambda_{p j}-l_{j}\right)-\frac{1}{2} \theta_{j}^{2} \tau^{2}\left(s_{0 j}\right)\right\}+o_{\mathrm{P}}(1)
$$

Similarly to the $\operatorname{Sig} \mathrm{D}$ case, the above $o_{\mathrm{P}}(1)$ term is uniform in $\left(\theta_{1}, \ldots, \theta_{k}\right)$ from any compact subset of $\mathbb{R}^{k}$ by construction.

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