Random Matrix Improved Covariance Estimation for a Large Class of Metrics

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Abstract
Relying on recent advances in statistical estimation of covariance distances based on random matrix theory, this article proposes an improved covariance and precision matrix estimation for a wide family of metrics. The method is shown to largely outperform the sample covariance matrix estimate and to compete with state-of-the-art methods, while at the same time being computationally simpler and faster. Applications to linear and quadratic discriminant analyses also show significant gains, therefore suggesting practical interest to statistical machine learning.

1. Introduction
Covariance and precision matrix estimation is a fundamental and simply posed, yet still largely considered, key problem of statistical data analysis, with countless applications in statistical inference. In machine learning, it is notably at the core of elementary methods as linear (LDA) and quadratic discriminant analysis (QDA) (McLachlan, 2004).

Estimation of the covariance matrix $C \in \mathbb{R}^{p \times p}$ based on $n$ independent (say zero mean) samples $x_1, \ldots, x_n \in \mathbb{R}^p$ is conventionally performed using the sample covariance matrix (SCM) $C \equiv \frac{1}{n} \sum_{i=1}^{n} x_i x_i^T$ (and its inverse using $\hat{C}^{-1}$). The estimate is, however, only consistent for $n \gg p$ and only invertible for $n \geq p$. Treating the important practical cases where $n \sim p$ and even $n \ll p$ has recently spurred a series of parallel lines of research. These directions rely either on structural constraints, such as “toeplitzification” procedures for Toeplitz covariance models (particularly convenient for time series) (Bickel et al., 2008; Wu & Pourahmadi, 2009; Vinogradova et al., 2015), or sparse constraints with LASSO and graphical LASSO-based approaches (Friedman et al., 2008) or, more interestingly for the present article, on exploiting the statistical independence in the entries of the vectors $x_i$.

(Ledoit & Wolf, 2004) proposes to linearly “shrink” $\hat{C}$ as $\hat{C}(\rho) \equiv \rho I_p + \sqrt{1-\rho^2} \hat{C}$ for $\rho > 0$ chosen to minimize the expected Frobenius distance $E[\|C - \hat{C}(\rho)\|_F]$ in the asymptotic $p, n \to \infty$ limit with $p/n \to c > 0$. Basic results from random matrix theory (RMT) are used here to estimate $\rho$ consistently. This procedure is simple and quite flexible and has been generalized in various directions (e.g., in (Couillet & McKay, 2014) with a robust statistics approach). However, the method only applies a naive homothetic map to each $\lambda_i(C)$ of $\hat{C}$ in order to better estimate $\lambda_i(C)$. A strong hope to recover a better approximation of the $\lambda_i(C)$’s then arose from (Silverstein & Bai, 1995; Silverstein & Choi, 1995) that provide a random matrix result relating directly the limiting eigenvalue distributions of $C$ and $\hat{C}$. Unfortunately, while estimating the $\lambda_i(\hat{C})$’s from the $\lambda_i(C)$’s is somewhat immediate, estimating the $\lambda_i(C)$’s backward from the $\lambda_i(\hat{C})$’s is a difficult task. (El Karoui et al., 2008) first proposed an optimization algorithm to numerically solve this problem, however, with little success as the method is quite unstable and has rarely been efficiently reproduced. (Mestre, 2008) later offered a powerful idea, based on contour integral, to consistently estimate linear functionals $\frac{1}{n} \sum_{i=1}^{n} f(\lambda_i(C))$ from the $\lambda_i(\hat{C})$’s. But $f$ is constrained to be very smooth (complex analytic) which prevents the estimation of the individual $\lambda_i(\hat{C})$’s. Recently, Ledoit and Wolf took over the work of El Karoui, which they engineered to obtain a more efficient numerical method, named QuEST (Ledoit & Wolf, 2015). Rather than inverting the Bai–Silverstein equations, the authors also proposed, with the same approach, to estimate the $\lambda_i(C)$’s by minimizing a Frobenius norm distance (Ledoit & Wolf, 2015) (named QuEST1 in the present article) or a Stein loss (Ledoit et al., 2018) (QuEST2 here).

These methods, although more stable than El Karoui’s initial approach, however, suffer several shortcomings: (i) they are still algorithmically involved as they rely on a series of fine-tuned optimization schemes, and (ii) they are only adaptable to few error metrics (Frobenius, Stein).

Inspired by Mestre’s approach and the recent work (Couil-
let et al., 2018), this article proposes a different procedure consisting in (i) writing \( C \) as the solution to \( \text{argmin}_{M>0} \delta(M, C) \) for a wide range of metrics \( \delta \) (Fisher, Batthacharyya, Stein’s loss, Wasserstein, etc.), (ii) based on (Couillet et al., 2018), using the fact that \( \delta(M, C) - \delta(M, X) \to 0 \) for some consistent estimator \( \hat{\delta} \), valid for all deterministic \( M \) and samples \( X = [x_1, \ldots, x_n] \in \mathbb{R}^{p \times n} \) having zero mean and covariance \( C \), and (iii) proceeding to a gradient descent on \( \hat{\delta} \) rather than on the unknown \( \delta \) itself. With appropriate adaptations, the estimation of \( C^{-1} \) is similarly proposed by solving instead \( \text{argmin}_{M>0} \delta(M, C^{-1}) \).

While only theoretically valid for matrices \( M \) independent of \( X \), the proposed method has several advantages: (i) it is easy to implement, technically simpler than QuEST and numerically fast, (ii) it is adaptable to a large family of distances and divergences, and, most importantly, (iii) simulations suggest that it systematically outperforms the SCM and is competitive with, and more robust than, QuEST.

The remainder of the article is organized as follows. Section 2 introduces preliminary notions and concepts on which are hinged our proposed algorithms, thereafter described in Section 3. Section 4 provides experimental validations and applications, including an improved version of LDA/QDA based on the proposed enhanced estimates.

**Reproducibility.** Matlab codes for the proposed estimation algorithms are available at https://github.com/maliktionoko/RMTCovEst and are based on Manopt, a Matlab toolbox for optimization on manifolds (Boumal et al., 2014).

### 2. Preliminaries

Let \( x_1, \ldots, x_n \in \mathbb{R}^p \) with \( x_i = C^{1/2} z_i \) for \( C \in \mathbb{R}^{p \times p} \) positive definite and \( z_1, \ldots, z_n \in \mathbb{R}^p \) independent random vectors of independent entries, where \( \text{E}[|z_i|] = 0 \) and \( \text{E}[|z_i|^2] = 1 \). We assume the following large dimensional regime for \( n \) and \( p \).

**Assumption 1** (Growth Rate). As \( n \to \infty \), \( p/n \to c \in (0, 1) \) and \( \lim \sup_p \max \{\|C^{-1}\|, \|C\|\} < \infty \) for \( \| \cdot \| \) the matrix operator norm.

Our objective is to estimate \( C \) and \( C^{-1} \) based on \( x_1, \ldots, x_n \) under the above large \( p, n \) regime. For simplicity of exposition and readability, we mostly focus on the estimation of \( C \) and more briefly discuss that of \( C^{-1} \).

Our approach relies on the following elementary idea:

\[
C \equiv \text{argmin}_{M>0} \delta(M, C)
\]

where, for some function \( f \),

\[
\delta(M, C) = \frac{1}{p} \sum_{i=1}^{p} f(\lambda_i(M^{-1}C))
\]

is a divergence (possibly a squared distance \( \delta = d^2 \) between the positive definite matrices \( M \) and \( C \), depending only on the eigenvalues of \( M^{-1}C \). Among divergences satisfying this condition, we find the *natural Riemannian distance* \( d^R_\delta \) (Bhatia, 2009), which corresponds to the Fisher metric for the multivariate normal distribution (Skovgaard, 1984); the *Bhatccharyya distance* \( d^B_\delta \) (Sra, 2013), which is close to the natural Riemannian distance while numerically cheaper; the *Kullback-Leibler divergence* \( d_{KL} \), linked to the likelihood and studied for example in (Moakher, 2012); the *Rényi divergence* \( d_{RE} \) for Gaussian \( x_i \)’s (Van Erven & Harremos, 2014); etc.\(^1\) Table 1 reports the explicit values of \( f \) for these divergences.

Since \( \delta(M, C) \) is not accessible as \( C \) is unknown, our approach exploits an estimator for \( \delta(M, C) \) which is consistent in the large \( n, p \) regime of Assumption 1.

To this end, our technical arguments are fundamentally based on random matrix theory, and notably rely on the so-called Stieltjes transform of eigenvalue distributions. For an arbitrary real-supported probability measure \( \theta \), the Stieltjes transform \( m_\theta : \mathbb{C} \setminus \text{supp}(\theta) \to \mathbb{C} \) is defined as

\[
m_\theta(z) = \int \frac{\theta(dt)}{t - z}.
\]

The key interest of the Stieltjes transform here lies in its relating the distributions of the eigenvalues of \( C \) and \( C \) as \( p, n \to \infty \) (Silverstein & Bai, 1995). Specifically, for arbitrary deterministic matrices \( M \), Stieltjes transform relations connect the empirical spectral (i.e., eigenvalue) distribution \( \nu_p \) of \( M^{-1/2}C \) to the empirical spectral distribution \( \mu_p \) of \( M^{-1/2}C \) (Couillet et al., 2018), defined as

\[
\mu_p = \frac{1}{p} \sum_{i=1}^{p} \delta_{\lambda_i(M^{-1}C)} \quad \text{and} \quad \nu_p = \frac{1}{p} \sum_{i=1}^{p} \delta_{\lambda_i(M^{-1}C)}.
\]

The connection goes as follows: first, from Cauchy’s integral formula \( f(t) = \frac{1}{2\pi i} \int_{\Gamma} f(z)(t - z)dz \) for \( \Gamma \) a complex contour enclosing \( t \), the metric \( \delta(M, C) \) in (1) relates to the Stieltjes transform \( m_\nu_p(z; M) \) through

\[
\delta(M, C) = \frac{1}{2\pi i} \int_{\Gamma} f(z)m_\nu_p(z; M)dz
\]

for \( \Gamma \subset \mathbb{C} \) a (positively oriented) contour surrounding the eigenvalues of \( M^{-1}C \). The notation \( m_\nu_p(z; M) \) reminds the dependence in \( M \) of \( m_\nu_p(z) \). Then, from the relation

\(^1\)The Frobenius distance does not fall into this setting but has already largely been investigated and optimized.
between the Stieltjes transforms \( m_{\nu_p} \) and \( m_{\mu_p} \), it follows from (Couillet et al., 2018) that, under Assumption 1, for all deterministic \( M \) of bounded operator norm,

\[
\delta(M, C) - \delta(M, X) \rightarrow 0 \tag{3}
\]

almost surely, where \( X = [x_1, \ldots, x_n] \) and

\[
\hat{\delta}(M, X) = \frac{1}{2\pi c f} \oint_{\Gamma} G(-m_{\hat{\mu}_p}(z; M)) \, dz \tag{4}
\]

with \( G \) such that \( G'(z) = g(z) = f'(1/z) \), \( \Gamma \) a contour surrounding the support of the almost surely limiting eigenvalue distribution of \( M^{-1} \hat{C} \) and \( \hat{\mu}_p = \frac{p}{n} \mu_p + (1 - \frac{p}{n})\delta_0 \) (and thus \( m_{\hat{\mu}_p}(z) = cm_{\nu_p}(z) + (1 - \frac{p}{n})z \)). Note that, by the linearity of \( G \) in (4), it is sufficient in practice to evaluate \( \hat{\delta}(M, X) \) for elementary functions (such as \( f(z) = z, f(z) = \log(z), f(z) = \log(1 + z) \)) etc. in order to cover most distances and metrics of interest (see again Table 1). Table 2 reports the values of \( G \) for such atomic functions \( f \).

Our main idea is to estimate \( C \) by minimizing the approximation \( \hat{\delta}(M, X) \) of \( \delta(M, C) \) over \( M \). However, it is important to note that, as discussed in (Couillet et al., 2018), the random quantity \( \hat{\delta}(M, X) \) may be negative with non-zero probability. As such, minimizing \( \hat{\delta}(M, X) \) over \( M \) may lead to negative solutions. Our proposed estimation method therefore consists in approximating \( C \) by the solution to the optimization problem

\[
\arg\min_{M \succ 0} h_X(M), \text{ with } h_X(M) \equiv (\hat{\delta}(M, X))^2. \tag{5}
\]

### 3. Methodology and Main Results

#### 3.1. Estimation Method

We solve (5) via a gradient descent algorithm on the Riemannian manifold \( S_n^{++} \) of positive definite \( n \times n \) matrices.

The Riemannian gradient \( \nabla h_X(M) \) of \( h_X \) at \( M \in S_n^{++} \) is defined via the directional derivative \( Dh_X(M)[\xi] \) of the functional \( h_X : S_n^{++} \rightarrow \mathbb{R}^+ \), at position \( M \in S_n^{++} \) and in the direction of \( \xi \in M_n \) (the vector space of symmetric \( n \times n \) matrices), by (Aisil et al., 2009)

\[
Dh_X(M)[\xi] = \langle \nabla h_X(M), \xi \rangle_{M_n}^{S_n^{++}}
\]

where \( \langle \cdot, \cdot \rangle_{S_n^{++}} \) is the Riemannian metric defined through

\[
\langle \eta, \xi \rangle_{M_n}^{S_n^{++}} = \text{tr} \left( M^{-1} \eta M^{-1} \xi \right).
\]

#### 3.2. Results

Differentiating \( \delta^2(M, X) \) at \( M \) in the direction \( \xi \) yields:

\[
Dh_X(M)[\xi] = -\frac{\delta(M, X)}{\pi c} \oint_{\Gamma} g(-m_{\hat{\mu}_p}(z, M)) Dm_{\hat{\mu}_p}(z, M) [\xi] \, dz.
\]

By using the fact that

\[
Dm_{\hat{\mu}_p}(z, M)[\xi] = \frac{c}{p} D \text{tr} \left( [M^{-1} \hat{C} - z I_p]^{-1} \right)[\xi] = \frac{c}{p} \text{tr} \left( [M^{-1} \hat{C} - z I_p]^{-2} M^{-1} \xi \right) = \langle \frac{c}{p} \text{sym} \left( \hat{C} [M^{-1} \hat{C} - z I_p]^{-2} \right), \xi \rangle_{M_n}^{S_n^{++}}
\]

where \( \text{sym}(A) = \frac{1}{2}(A + A^T) \) is the symmetric part of \( A \in \mathbb{R}^{p \times p} \), we retrieve the gradient of \( h_X(M) \) as

\[
-\frac{\pi c p}{\delta(M, X)} \nabla h_X(M) = \oint_{\Gamma} g(-m_{\hat{\mu}_p}(z; M)) \text{sym} \left( \hat{C} [M^{-1} \hat{C} - z I_p]^{-2} \right) \, dz
\]

(6)

(recall that the right-hand side still depends on \( X \) implicitly through \( \hat{\mu}_p \) and \( \hat{C} \)).

Once \( \nabla h_X \) estimated, every gradient descent step in \( S_n^{++} \) corresponds to a small displacement on the geodesic starting at \( M \) and towards \( -\nabla h_X(M) \), defined as the curve

\[
\mathbb{R}_+ \rightarrow S_n^{++} \quad t \mapsto M^t \exp \left( -t M^{-\frac{1}{2}} \nabla h_X(M) M^{-\frac{1}{2}} \right) M^{\frac{1}{2}}
\]

where, for \( A = U \Lambda U^T \in S_n^{++} \) in its spectral decomposition, \( \exp(A) \equiv U \exp(\Lambda) U^T \) (with \( \exp \) understood here applied entry-wise on the diagonal elements of \( \Lambda \)).

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<table>
<thead>
<tr>
<th>Divergences</th>
<th>( f(z) )</th>
<th>( G(z) )</th>
</tr>
</thead>
<tbody>
<tr>
<td>( d_R^2 )</td>
<td>( \log^2(z) )</td>
<td>( z (\log^2(z) - 2 \log(z) + 2) )</td>
</tr>
<tr>
<td>( d_B^2 )</td>
<td>( -\frac{1}{q} \log(z) + \frac{1}{2} \log(1 + z) - \frac{1}{2} \log(2) )</td>
<td>( \log(z) - z \log(z) + z )</td>
</tr>
<tr>
<td>( \delta_{KL} )</td>
<td>( \frac{1}{2} z - \frac{1}{q} \log(z) - \frac{1}{2} )</td>
<td>( \log(2) )</td>
</tr>
<tr>
<td>( \delta_{\alpha R} )</td>
<td>( \frac{2}{\alpha - 1} \log(\alpha + (1 - \alpha)z) + \frac{1}{2} \log(z) )</td>
<td>( \log(z) )</td>
</tr>
</tbody>
</table>

Table 1. Distances \( d \) and divergences \( \delta \), and corresponding \( f(z) \).

Table 2. \( G(z) \) and \( F(z) \) for “atomic” \( f(z) \) functions used in most distances and divergences under study; here \( s > 0 \) and \( z \in \mathbb{C} \).
That is, letting $M_0, M_1, \ldots$ and $t_0, t_1, \ldots$ be the successive iterates and step sizes of the gradient descent, we have, for some given initialization $M_0 \in S_\mathcal{n}^{++}$,
\begin{equation}
M_{k+1} = M_k^\frac{1}{2} \exp \left( -t_k M_k^{-\frac{1}{2}} \nabla h_X(M_k) M_k^{-\frac{1}{2}} \right) M_k^\frac{1}{2}.
\end{equation}
(7)

Our proposed method is summarized as Algorithm 1.

Algorithm 1 Proposed estimation algorithm.

Require $M_0 \in S_\mathcal{n}^{++}$.

Repeat $M \leftarrow M^\frac{1}{2} \exp \left( -t M^{-\frac{1}{2}} \nabla h_X(M) M^{-\frac{1}{2}} \right) M^\frac{1}{2}$ with $t$ either fixed or optimized by backtracking line search.

Until Convergence.

Return $M$.

We conclude this section by an important remark on the fundamental limitations of the proposed algorithm.

Remark 1 (Approximation of $\delta(M_k, C)$ by $\hat{\delta}(M_k, X)$). It is fundamental to understand the result from (Couillet et al., 2018) at the heart of the proposed method. There, it is precisely shown that, for every deterministic sequence of matrices $\{M^{(p)}\}$ and $\{C^{(p)}\}$ with $M^{(p)}, C^{(p)} \in \mathbb{R}^{p \times p}$ and $\max \{\|C^{(p)}\|, \|M^{(p)}\|\} < K$ for some constant $K$ independent of $p$, we have that, for $X^{(p)} = [x_1^{(p)}, \ldots, x_n^{(p)}]$ with $x_i^{(p)} = C^{(p)} \frac{1}{2} z_i^{(p)}$ and $z_i^{(p)}$, i.i.d. vectors of i.i.d. zero mean and unit variance entries,
\begin{equation}
\delta(M^{(p)}, C^{(p)}) - \hat{\delta}(M^{(p)}, X^{(p)}) \rightarrow 0
\end{equation}
almost surely as $n, p \rightarrow \infty$ and $p/n \rightarrow c \in (0, 1)$. This result seems to suggest that $\hat{\delta}(M_k, X)$ in our algorithm is a good approximation for the sought for $\delta(M_k, C)$. This, however, only holds true so long that $M_k$ is independent of $X$ which clearly does not stand when proceeding to successive gradient descent steps in the direction of $\nabla h_X(M)$ which depends explicitly on $X$. As such, while initializations with, say, $M_0 = I_p$, allow for a close approximation of $\delta(M_0, C)$ in the very first steps of the descent, for larger values of $k$, the descent is likely to drive the optimization in less accurate directions.

Remark 1 is in fact not surprising. Indeed, finding the minimum of $\delta(M, C)$ over $M \succ 0$ would result in finding $C$, which cannot be achieved for unconstrained matrices $\mathcal{C}$ and for non vanishing values of $p/n$. Figure 1 provides a typical evolution of the distance $\delta(M_k, C)$ versus its approximation $\hat{\delta}(M_k, X)$ at the successive steps $k = 1, 2, \ldots$ of Algorithm 1, initialized at $M_0 = I_p$. As expected, the difference $|\hat{\delta}(M_k, X) - \delta(M_k, C)|$, initially small (at $k = 1$, $\delta(I_p, X) \cong \delta(I_p, C)$), increases with $k$, until the gradient vanishes and the divergence $\delta(M_k, C)$ converges.

3.2. Practical Implementation

In order to best capture the essence of Algorithm 1, as well as its various directions of simplification and practical fast implementation, a set of important remarks are in order.

First, the generic computation of $M_{k+1}$ in Equation (7) may be numerically costly, unless $M_k$ and $\nabla h_X(M)$ are co-diagonalizable by an orthogonal matrix. In this case, for $M_k = U\Omega_k U^T$ and $\nabla X(h(M_k)) = U\Delta_k U^T$, we have the recursion
\begin{equation}
[\omega_{k+1}]_i = [\omega_k]_i \exp \left( -t \frac{\delta_k}{|\omega_k|} \right)
\end{equation}
where $\delta_k = \text{diag}(\Delta_k)$ and $\omega_k = \text{diag}((\Omega_k))$. In particular, if $M_0$ has the same eigenvectors $\hat{U} \in \mathbb{R}^{p \times p}$ as $\mathcal{C}$ (for instance, $M_0 = \alpha I_n + \sqrt{1 - \alpha^2} \mathcal{C}$, a linear shrinkage of $\mathcal{C}$ for some $\alpha \in [0, 1]$), then for $k \geq 0$, $M_k = U\Omega_k U^T$, with $\Omega_k$ recursively defined through (8). Thus, for this $M_0$, the ultimate estimator of $\mathcal{C}$ shares the same eigenvectors as $\hat{\mathcal{C}}$ and reduces to a “non-linear shrinkage” procedure, similar to (Ledoit & Wolf, 2015).

Remark 2 (On initialization with linear shrinkage). Equation (6) with $M \propto \mathcal{C}$ shows that $\nabla h_X(\mathcal{C}) \propto \mathcal{C}$; thus, initialized at $M_0 = \hat{\mathcal{C}}$ (a particularly inappropriate choice since $\hat{\mathcal{C}}$ strongly depends on $X$), the iterations $M_k$ live on the line $\{\beta \hat{\mathcal{C}}, \beta \in \mathbb{R}\}$. This situation must be avoided. On the opposite, initializing at $M_0 = I_p$ ensures $\hat{d}(M_0, X) \approx d(M_0, C)$ although here $M_0$ starts far off the solution. The aforementioned linear shrinkage initialization $M_0 = \alpha I_n + \sqrt{1 - \alpha^2} \mathcal{C}$ is an appropriate tradeoff: for small $n/p$, $M_0 \approx I_p$ ensures $\hat{d}(M_0, X) \approx d(M_0, C)$, hence a fast and trustable initial gradient step progression; for large $n/p$, $M_0 \approx \hat{\mathcal{C}} \approx \mathcal{C}$ is almost deterministic, and thus again $\hat{d}(M_0, X) \approx d(M_0, C)$.

Remark 3 (On the eigenvectors of $M_\infty$). We have seen above that $\mathcal{H} = \{\hat{U}D\hat{U}^T \mid D\text{ diagonal}\}$ is a stable set of Algorithm 1 in the sense that $M_k \in \mathcal{H} \Rightarrow M_{k+1} \in \mathcal{H}$. One may thus wonder if $\mathcal{H}$ is also a global attractor: i.e., does every trajectory $\{M_1, M_2, \ldots\}$ necessarily converge to $\mathcal{H}$? Extensive simulations initialized randomly (say with $M_0$ a random Wishart matrix) indeed suggest that, after a few iterates, the eigenvectors of $M_k$ do converge to those of $\mathcal{C}$. This is, however, not everywhere true. Indeed, in the extreme scenario where $M_0 = \mathcal{C}, 0 = \delta(C, C) \approx \delta(C, X)$ which consistently estimates zero. For large $n, p$ (and irrespective of $n/p$), the gradient descent does not progress much from $M_0$ and is thus unlikely converging within $\mathcal{H}$ (which would mean the existence of a $D$ such that $|\delta(U\hat{D}\hat{U}^T, X)| < |\delta(C, X)|$).

A further direction of simplification of Algorithm 1 relates to the fact that, for generic values of $M_0$ (notably having eigenvectors different from those of $\hat{C}$), Equation (7)
is computationally expensive to evaluate. A second-order simplification for small $t$ is often used in practice (Jeuris et al., 2012), as follows

$$M_{k+1} = M_k - t \nabla h_X(M_k)$$

$$+ \frac{t^2}{2} \nabla h_X(M_k) M_k^{-1} \nabla h_X(M_k) + O(t^3).$$

Simulations with this approximation suggest almost no difference in either the number of steps until convergence or accuracy of the solution.

### 3.3. Estimation of $C^{-1}$

In our framework, estimating $C^{-1}$ rather than $C$ can be performed by minimizing $\delta(M, C^{-1})$ instead of $\delta(M, C)$. In this case, under Assumption 1, (3) now becomes

$$\delta(M, C^{-1}) - \delta^{\text{inv}}(M, X) \to 0$$

almost surely, for every deterministic $M$ of bounded operator norm and $X = [x_1, \ldots, x_n]$, where

$$\delta^{\text{inv}}(M, X) \equiv \frac{1}{2 \pi i c} \oint \mathcal{F} \left( -m_p^{\text{inv}}(z; M) \right) dz$$

for $\mathcal{F}$ such that $\mathcal{F}'(z) \equiv f(z)$, $\Gamma$ a contour surrounding the support of the almost sure limiting eigenvalue distribution of $M \tilde{C}$ and $m_p^{\text{inv}} = \frac{1}{p} m_p^{\text{inv}} + (1 - \frac{1}{p}) \delta_0$, where $m_p^{\text{inv}} \equiv \frac{1}{p} \sum_{i=1}^p \delta_{\lambda_i(M \tilde{C})}$. The cost function to minimize under this setting is now given by $h_X(M) \equiv (\delta^{\text{inv}}(M, X))^2$ with gradient $\nabla h_X^{\text{cov}}(M)$ satisfying

$$ivp \nabla h_X^{\text{cov}}(M, X) =$$

$$\oint \mathcal{F} \left( -m_p^{\text{inv}}(z; M) \right) \text{sym} \left( M \tilde{C} - zI_p \right)^{-2} dz.$$

With these amendments, Algorithm 1 can be adapted to the estimation of $C^{-1}$. Table 2 provides the values of $F$ for the atomic functions $f$ of interest.

### 3.4. Application to Explicit Metrics

Algorithm 1 is very versatile as it merely consists in a gradient descent method for various metrics $f$ through adaptable definitions of the function $h_X(M) = \delta(M, X)^2$ and its resulting gradient. Yet, because of the integral form assumed by the gradient (Equation (6)), a possibly computationally involved complex integration needs to be numerically performed at each gradient descent step.

In this section, we specify closed-form expressions for the gradient for the atomic $f$ functions of Table 2 (which is enough to cover the list of divergences in Table 1).

#### 3.4.1. Estimation of $C$

Let us denote

$$\nabla h_X(M) \equiv 2 \delta(M, X) \cdot \text{sym} \left( \tilde{C} \cdot V \Lambda \psi V^{-1} \right)$$

where $V$ are the eigenvectors of $M^{-1} \tilde{C}$ and $\Lambda \psi$ is to be determined for each $f$.

For readability in the following, let us denote $\lambda_i \equiv \lambda_i(M^{-1} \tilde{C})$, $i \in \{1, \ldots, p\}$, the eigenvalues of the matrix $M^{-1} \tilde{C}$ and $\xi_1, \ldots, \xi_p$ the eigenvalues of $\Lambda = \frac{1}{\sqrt{\lambda}} \Lambda^\top$ with $\Lambda = \text{diag}(\lambda_1, \ldots, \lambda_p)$ and $\lambda = (\lambda_1, \ldots, \lambda_p)^\top$. Finally, for $s > 0$, let $\kappa_s \in (-1/(s(1-p/n)), 0)$ be the unique negative number $t$ solution of the equation (see Couillet et al., 2018 for details) $m_{\tilde{p}}(t) = -s$.

With these notations at hand, following the derivations in (Couillet et al., 2018), we have the following determinations for $\Lambda \psi$.

**Proposition 1** (Case $f(t) = t$).

$$[\Lambda \psi]_{kk} = -\frac{1}{c} + \frac{1}{p} \sum_{i=1}^p \frac{1}{m_p'(\xi_i)} \left( \lambda_k - \xi_i \right)^2$$

with $m_p'$ the derivative of $m_p$.

**Proposition 2** (Case $f(t) = \log(t)$). $[\Lambda \psi]_{kk} = \frac{-1}{p \lambda_k}$.

**Proposition 3** (Case $f(t) = \log(1 + st)$). For $s > 0$,

$$[\Lambda \psi]_{kk} = \frac{-1}{p(\lambda_k - \kappa_s)}$$

**Proposition 4** (Case $f(t) = \log^2(t)$). For $f(t) = \log^2(t)$,

$$[\Lambda \psi]_{kk} = \frac{2}{p} \log(\lambda_k) \left[ \sum_{i=1}^p \frac{1}{\lambda_k - \xi_i} - \sum_{i \neq k}^p \frac{1}{\lambda_k - \lambda_i} - \frac{1}{\lambda_k} \right]$$

$$- \frac{2}{p} \sum_{i=1}^p \log(\xi_i) + \frac{2}{p} \sum_{i \neq k}^p \log(\lambda_i) - \frac{2 - 2 \log(1 - c)}{p \lambda_k}.$$
These results unfold from residue calculus for entire functions \( f \) or advanced complex integration tools for logarithmic functions (see supplementary material). Linear combinations of these formulas provide analytic expressions for the gradients of all aforementioned divergences.

Besides, recall (Section 3.2) that, if \( M_0 \) and \( \hat{C} \) have the same eigenvectors, then so do \( M_k \) and thus \( \nabla h_X(M_k) \) for all \( k \geq 0 \). In this case, the eigenvalues \( \lambda_i = [\hat{\lambda}_k]_i \) of \( M_k^{-1} \hat{C} \) in the above propositions are simply given by \( [\hat{\lambda}_k]_i = \frac{[\omega_i]}{[\tilde{\omega}_i]} \) with \( [\omega_i]_i = \lambda_i(M_k) \) and \( \ell_i = \lambda_i(\hat{C}) \). These considerations combined yield Algorithm 2.

**Algorithm 2** Improved estimation from linear shrinkage.

**Require** \( \omega_0 = \alpha + \sqrt{1 - \alpha^2} \ell_i \) for \( \hat{C} = U \text{diag}(\ell_i) U^T \).

**Repeat** for each \( i \)

**Step 1:** \( [\lambda_k]_i = \frac{[\omega_i]}{[\tilde{\omega}_i]} \)

**Step 2:** \( [\delta_k]_i = h([\lambda_k]_i) \) with \( h \) obtained from Prop. 1–4.

**Step 3:** \( [\omega_k+1]_i = [\omega_k]_i \exp \left( -t [\delta_k]_i \right) \), \( t > 0 \)

**Until** Convergence at step \( K \).

**Return** \( M_K = \hat{U} \Omega_K \hat{U}^T \), with \( \Omega_K = \text{diag}([\omega_K]_i) \).

### 3.4.2. ESTIMATION OF \( C^{-1} \)

Similarly, for the problem of estimating \( C^{-1} \), recalling Section 3.3, we may denote

\[
\nabla h_X^{\text{inv}}(M) \equiv 2\delta^{\text{inv}}(M, X) \cdot \text{sym} \left( M : V_{\text{inv}} \Lambda^{\text{inv}} V_{\text{inv}}^{-1} \right)
\]

with \( V_{\text{inv}} \) the eigenvectors of \( M \hat{C} \). We redefine \( \lambda_i \equiv \lambda_i(M \hat{C}) \), and \( \xi_1, \ldots, \xi_p \) the eigenvalues of \( \Lambda = \frac{1}{n} \sqrt{\Lambda \Lambda^T} \) with \( \Lambda = \text{diag}(\lambda_1, \ldots, \lambda_p) \) and \( \lambda = (\lambda_1, \ldots, \lambda_p)^T \).

Again, for \( s > 0 \), let \( \kappa_s < 0 \) be the only negative real number \( t \) solution of \( m(\Lambda)_{\text{inv}}(t) = -\frac{1}{2} \). With the same approach as above, we obtain the following values for \( \Lambda^{\text{inv}} \).

**Proposition 5** (Case \( f(t) = t \)). \( [\Lambda^{\text{inv}}]_{kk} = -\frac{1}{\lambda_k} \).

**Proposition 6** (Case \( f(t) = \log(t) \)). \( [\Lambda^{\text{inv}}]_{kk} = -1 \).

**Proposition 7** (Case \( f(t) = \log(1 + st) \)). For \( s > 0 \),

\[
[\Lambda^{\text{inv}}]_{kk} = \frac{\lambda_k}{\lambda_k - \kappa_s} - 1.
\]

**Proposition 8** (Case \( f(t) = \log^2(t) \)). For \( f(t) = \log^2(t) \),

\[
[\Lambda^{\text{inv}}]_{kk} = -\frac{2}{p} \log(\lambda_k) \left[ \sum_{i=1}^{p} \frac{\lambda_k}{\lambda_k - \xi_i} - \sum_{i=1}^{p} \frac{\lambda_k}{\lambda_k - \lambda_i} - 1 \right]
\]

\[
+ \frac{2}{p} \sum_{i=1}^{p} \frac{\lambda_k \log(\xi_i)}{\lambda_k - \xi_i} - \frac{2}{p} \sum_{i=1}^{p} \frac{\lambda_k \log(\lambda_i)}{\lambda_k - \lambda_i} + \frac{2}{p} \log(1 - c).
\]

### 4. Experimental Results

This section introduces experimental results on the direct application of our proposed method to the estimation of \( C \) and \( C^{-1} \) as well as on its use as a plug-in estimator in more advanced procedures, here in the scope of linear and quadratic discriminant analyses (LDA/QDA).

#### 4.1. Validation on synthetic data

This first section provides simulations on the estimation of \( C \) and \( C^{-1} \) based on the Fisher distance for several matrices \( C \).

The Fisher distance, being the “natural” Riemannian distance to compare covariance matrices in \( S_{++}^n \), better agrees with the proposed estimation strategy through gradient descents in \( S_{++}^n \). Besides, Theorem 4 in (Smith, 2005) establishes an exact formula for the Cramer–Rao bound (CRB) on unbiased estimators of \( C \); although the compared estimators of \( C \) are likely all biased and that \( M_0 \) initializations may by chance bring additional information disrupting a formally fair CRB comparison, the CRB at least provides an indicator of relevance of the estimators.

The chosen matrix \( C \) in the following is either:

1. [Wishart] a random \((p \times p)\) standard Wishart matrix with \( 2p \) degrees of freedom,
2. [Toeplitz] a matrix defined by \( C_{ij} = a^{n+j-i} \),
3. [Discrete] a matrix with random uniform eigenvectors and eigenvalues \( \{1, 2, 3, 4\} \) each with multiplicity \( p/4 \).

Figure 2 (for the estimation of \( C \)) and Figure 3 (for \( C^{-1} \)) report comparative performances for the SCM, QuEST1, QuEST2, and our proposed estimator, with the latter three initialized at \( M_0 \) the linear shrinkage estimate of (Ledoit & Wolf, 2004). In the figures, “SCM th” refers to the analytical approximation of \( \delta(C, \hat{C}) \) defined in Remark 4 below. Both our proposed method and QuESTx perform competitively and systematically better than the sample covariance matrix (and sometimes significantly overtake the CRB). Yet, QuESTx have a larger variance, which is in part due to the occurrence of outlying unstable solutions (which our approach does not seem to suffer). This unstable behavior can be traced back to QuEST being founded on a computationally intense and involved numerical procedure composed of six intricate steps of finely-tuned optimization schemes (see implementation details in (Ledoit & Wolf, 2017)) aiming at stabilizing El Karoui’s initial approach (El Karoui et al., 2008), where our approach is a much more stable gradient descent method.

Comparing QuESTx and our method is a difficult task though. Both rely on different “approximations”: QuESTx...
use numerical methods to solve several non-convex inverse problems; our technique is limited by the dependence on $X$ of the gradient steps. Each approximation loses accuracy under various conditions, making theoretical comparison challenging. The numerical approach of QuESTx also makes the algorithm both more intricate and slower. Table 3 compares the typical algorithm duration.

**Remark 4** (Consistent estimator for $\delta(C, \tilde{C})$). With the same technical tools from (Couillet et al., 2018), it is straightforward to estimate the distance $\delta(C, \tilde{C})$. Indeed, $\delta(C, \tilde{C}) = \frac{1}{2\pi^2} \int f(z)m_a(z)dz$ for $m_a(z)$ the Stieltjes transform of the eigenvalue distribution of $C^{-1}\tilde{C}$; the limiting distribution of the latter is the popular Marcenko-Pastur law (Marčenko & Pastur, 1967), the expression of which is well known. The estimate is denoted “SCM th” in Figures 2–3. The observed perfect match between limiting theory and practice confirms the consistency of the random matrix approach even for not too large $p, n$.

### 4.2. Application to LDA/QDA

As pointed out in the introduction, the estimation of the covariance and inverse covariance matrices of random vectors are at the core of a wide range of applications in statistics, machine learning and signal processing. As a basic illustrative example, we focus here on linear discriminant analysis (LDA) and quadratic discriminant analysis (QDA). Both exploit estimates covariance matrices of the data or their inverse in order to perform the classification.

Suppose $x_1^{(1)} \ldots x_{n_1}^{(1)} \sim N(\mu_1, C_1)$ and $x_1^{(2)} \ldots x_{n_2}^{(2)} \sim N(\mu_2, C_2)$ are two sets of random independent $p$-dimensional training vectors forming two classes of a Gaussian mixture. The objective of LDA and QDA is to estimate the probability for an arbitrary random vector $x$ to belong to either class by replacing the genuine means $\mu_a$ and covariances $C_a$ by sample estimates, with in the case of LDA the underlying (possibly erroneous) assumption that $C_1 = C_2$. Defining $C \equiv \frac{n_1}{n_1 + n_2} C_1 + \frac{n_2}{n_1 + n_2} C_2$, the classification rules for LDA and QDA for data point $x$ depend on the signs of the respective quantities:

$$
\delta_{x}^{\text{LDA}} = (\hat{\mu}_1 - \hat{\mu}_2)^T \tilde{C}^{-1} x + \frac{1}{2} \hat{\mu}_2^T \tilde{C}^{-1} \mu_2 - \frac{1}{2} \hat{\mu}_1^T \tilde{C}^{-1} \mu_1
$$

$$
\delta_{x}^{\text{QDA}} = \frac{1}{2} x^T (\hat{C}_1^{-1} - \tilde{C}_1^{-1}) x + (\hat{\mu}_1^T \hat{C}_1^{-1} - \hat{\mu}_2^T \tilde{C}_2^{-1}) x + \frac{1}{2} \hat{\mu}_1^T \tilde{C}_1^{-1} \mu_1 + \frac{1}{2} \log \det \tilde{C}_1^{-1} - \log \frac{n_2}{n_1}
$$

where $\hat{\mu}_a = \frac{1}{n_a} \sum_{i=1}^{n_a} x_i^{(a)}$ is the sample estimate of $\mu_a$ and $\tilde{C}_a^{-1}$ are some estimate of $C_a^{-1}$ while $\hat{C}_a \equiv \frac{n_1}{n_1 + n_2} C_1 + \frac{n_2}{n_1 + n_2} \tilde{C}_2$ with $\tilde{C}_a$ the estimation of $C_a$. As such, here, LDA only exploits estimations of $C_1$ and $C_2$ (before inverting the estimated average), while QDA focuses on estimating the inverses $C_1^{-1}$ and $C_2^{-1}$.
The first three displays in Figures 4 and 5 compare the accuracy of LDA/QDA for $C_1$ and $C_2$ chosen among Wishart and Toeplitz matrices, and $\mu_2 = \mu_1 + \frac{50}{p}$ for LDA, $\mu_2 = \mu_1 + \frac{1}{p}$ for QDA (to avoid trivial classification). The bottom right displays are applications to EEG data from (Andrzejak et al., 2001). The dataset contains 5 subsets: sets A and B collected from healthy, C, D, E from epileptic subjects. The graph presents all combinations of binary classes between healthy and epileptic subjects (e.g., A/E for subsets A and E). For most settings, our proposed algorithm almost systematically outperforms competing methods, with QuEST1 and QuEST2 exhibiting a much less stable behavior and quite weak performances in all synthetic scenarios.

5. Discussion and Concluding Remarks

We have proposed in this work a random matrix framework for the estimation of covariance and precision matrices. Unlike alternative state-of-the-art techniques that attempt to invert the fundamental Bai–Silverstein equations (Silverstein & Bai, 1995), our method relies on basic gradient descent in $S_+^n$ that, in addition to performing competitively (if not better), is computationally simpler.

While proposed here for metrics involving the eigenvalues of products of covariance matrices, the method flexibly adapts to further divergences solely depending on eigenvalues. The same framework can notably be applied to the Wasserstein distance between zero-mean Gaussians. A reservation nonetheless remains on the need for $n > p$: (Couillet et al., 2018) shows that this problem can be partially avoided for some divergences (not for the Fisher distance though); some workarounds are being investigated.

Our approach also suffers a profound limitation: $\delta(M, X)$ only estimates $\delta(M, C)$ for $M$ independent of $X$. This poses a formal problem when implemented in a gradient descent. This needs be tackled: (i) either by estimating the introduced bias so to infer the loss incurred or, better, (ii) by accounting for the dependence to provide a further estimator $\hat{\delta}(M, X, C)$ of $\delta(M, X, C)$ for all $X$-dependent matrices $M(X)$ following a specific form. Notably, since $\mathcal{H} = \{ U D U^T \}$, with $U$ the eigenvectors of $C$, is stable through gradient descent (see Remark 3), a first improvement consists in estimating consistently $\delta(U D U^T, C)$ for deterministic diagonal matrices $D$.

We conclude by emphasizing that modern large dimensional statistics have lately realized that substituting large covariance matrices by their sample estimators (or even by improved covariance estimators) is in general a weak approach, and that one should rather focus on estimating some ultimate functional (e.g., the result of a statistical test) involving the covariance (see, e.g., (Mestre & Lagunas, 2008) in array processing or (Yang et al., 2015) in statistical finance). It is to be noted that our proposed approach is consistent with these considerations as various functionals of $C$ can be obtained from Equation (2), from which similar derivations can be performed.
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References


