

Testing distributional assumptions using a continuum of moments*

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Abstract

We propose specification tests for parametric distributions that compare the potentially complex theoretical and empirical characteristic functions using the continuum of moment conditions analogue to an overidentifying restrictions test, which takes into account the correlation between influence functions for different argument values. We derive its asymptotic distribution for fixed regularization parameter and when this vanishes with the sample size. We show its consistency against any deviation from the null, study its local power and compare it with existing tests. An extensive Monte Carlo exercise confirms that our proposed tests display good power in finite samples against a variety of alternatives.

Keywords: Characteristic function, Complex Gaussian process, Consistent tests, Continuum of moment conditions, Goodness-of-fit, GMM, Tikhonov regularization.

JEL: C01, C12, C52.

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1 Introduction

Goodness-of-fit tests are important to assess whether a parametric distribution provides an appropriate representation of the data. These tests can be divided in two main categories: (i) directional tests, which are designed to have power against specific alternatives, such as Neyman smooth test (see Neyman, 1937 and Rayner and Best, 1989), Jarque and Bera's (1980) test of normality, as well as those proposed by Sefton (1992), Fiorentini, Sentana and Calzolari (2003), Bontemps and Meddahi (2005, 2012), Mencía and Sentana (2012) and Tuvaandorj and Zinde-Walsh (2014) among many others; (ii) omnibus tests, which are consistent against any alternative to the null hypothesis, for instance the integrated conditional moment test of Bierens (1982) and Bierens and Ploberger (1997), the conditional Kolmogorov test of Andrews (1997), and the copula goodness-of-fit test of Genest, Huang and Dufour (2013). Our proposed tests fall in this second category.

In particular, our testing procedure is based on the difference between the empirical and theoretical characteristic functions (CF) for all possible values of their argument. This gives rise to a continuum of moments in a L^2 space. Our aim is to construct a J test for overidentifying restrictions based on these moments, as in Hansen (1982). However, what plays the role of the covariance matrix in his test becomes now a covariance operator, whose inverse is unbounded. Therefore, some regularization is needed to stabilize the inverse. We propose to use Tikhonov regularization (see Kress, 1999) and consider two types of tests. The first one uses a fixed value of the regularization parameter α . Given that α can be regarded as a bandwidth, this approach is analogous to the fixed b asymptotics used in Kiefer and Vogelsang (2002). The second type of tests allows α to converge to zero at an appropriate rate, in which case our proposed test is closer in spirit to Hansen (1982)'s J test. In this second instance, however, the statistics would tend to a diverging χ^2 with infinite degrees of freedom. For that reason, we center and rescale it following the procedure put forward by Carrasco and Florens (2000), who presented this type of test for the first time. Note that Carrasco and Florens (2000) assume that the moment conditions are real whereas here we work with complex moment conditions.

We will consider various versions of our proposed tests depending on whether the parameter vector θ is known in advance or replaced by a consistent estimator, and whether we make use of the analytical expression for the covariance operator or estimate it. We will derive the asymptotic distribution of our tests under the null hypothesis and under local alternatives. We will also characterize the alternatives for which our tests have maximum power.

The advantages of using the CF are multiple: (a) in some important examples, the distribution function is only known in integral form whereas the CF has a closed form expression,

as in the cases of stable distributions and affine diffusions (see Singleton (2001) and Carrasco, Chernov, Florens, and Ghysels (2007)); (b) handling multivariate random variables can be done just as easily as the scalar case; (c) our tests have the same form and are computed in the same manner for any CF tested; (d) our tests are consistent against any alternative to the null hypothesis.

The CF raises specific challenges as the CF-based moment condition converges to a complex Gaussian process but some results valid for real processes are not directly applicable to complex processes. To derive the asymptotic distributions of our tests, we adapt some results recently developed by Ducharme, de Micheaux, and Marchina (2016) for complex random vectors to complex processes. For results on the weak convergence of the empirical CF, see Csörgö (1981) and Wells (1992). Various tests based on the empirical CF have been previously proposed: Feuerverger and Mureika (1977), Epps and Pulley (1983), Hall and Welsh (1983), Baringhaus and Henze (1988), Ghosh and Ruymgaart (1992), Fan (1997), Hong (1999), Su and White (2007), Chen and Hong (2010), and Leucht (2012) among others. The most closely related paper is that of Bierens and Wang (2012), which focuses on tests for parametric conditional distributions. Recently, Bierens and Wang (2017) extended their tests to time-series. The main difference with ours is that we “weight” the continuum of moment conditions by the inverse of their covariance operator. Our work is also related to Dufour and Valery (2016), who propose a regularized Wald test to deal with the singularity of the covariance matrix.

The remainder of the paper is organized as follows. We introduce our tests in Section 2 and derive the asymptotic properties of the J test with fixed regularization parameter α and known (unknown) θ in Section 3 (4). Next, we study the J test with vanishing α in Section 5. Finally, Section 6 presents the results of our Monte Carlo simulations while Section 7 concludes. All the proofs are collected in the appendix and computational aspects as well as additional figures are included in the online Supplemental Appendix.

2 Presentation of the tests and overview

Assume we observe a sample of random variables X_1, X_2, \dots, X_n independent and identically distributed (*iid*) taking their values on \mathbb{R}^q with $q \geq 1$. The X_j have probability density function (pdf) $f(x; \theta)$ indexed by a finite dimensional parameter θ , which may be known or unknown, and CF $\psi(t; \theta) = E[e^{itX}]$, where $t \in \mathbb{R}^q$ is its argument. As is well known, $f(x; \theta)$ and $\psi(t; \theta)$ are intimately related because the former is the Fourier transform of the latter, i.e.

$$\psi(t; \theta) = \int e^{itx} f(x; \theta) dx. \quad (1)$$

Figure B1 in the Supplemental Appendix presents the CFs for the univariate distributions that we consider in our Monte Carlo study, namely, a standard normal, as well as standardized (zero mean - unit variance) versions of the uniform and $\chi^2(2)$ distributions, and a Cauchy distribution with location and scale 0 and 1, respectively. Given that the first two examples and the Cauchy are symmetrically distributed around 0, the CF is real and symmetric around 0. In contrast, it contains an (odd) imaginary component in the case of the asymmetric chi-square.

We are interested in testing $H_0 : \psi = \psi_0(\cdot; \theta_0)$, where ψ_0 is a known CF and θ_0 is some element of $\Theta \subset \mathbb{R}^p$. Our testing procedures are based on the difference between the empirical and theoretical CFs. Specifically, the relevant influence functions are

$$\hat{h}(t; \theta) = \frac{1}{n} \sum_{j=1}^n h_j(t; \theta), \quad (2)$$

$$h_j(t; \theta) = e^{itX_j} - \psi_0(t; \theta). \quad (3)$$

This gives rise to a continuum of moments since under the null $E[h_j(t; \theta_0)] = 0$ for all $t \in \mathbb{R}^q$.

Let π be a probability density function with support \mathbb{R}^q . Then, the function $h_j(t; \theta)$ is a random element of $L^2(\pi)$, the space of complex-valued functions which are square integrable with respect to the density π . The inner product on this space is defined for any functions f and g of $L^2(\pi)$ as $\langle f, g \rangle = \int f(t) \overline{g(t)} \pi(t) dt$, where the bar denotes the complex conjugate. $L^2(\pi)$ is a Hilbert space and we will work on this space to derive the asymptotic distribution of our test statistics.

By the central limit theorem of *iid* random elements of a separable Hilbert space (see e.g. proof of Theorem 9 in Rackauskas and Suquet, 2006), we have that under H_0 , as n goes to infinity

$$\sqrt{n} \hat{h}(\cdot; \theta_0) \Rightarrow \mathcal{CN}(0, K, R)$$

in $L^2(\pi)$, where $\mathcal{CN}(0, K, R)$ denotes a complex Gaussian process of $L^2(\pi)$. This process is characterized by its mean, its covariance operator K , which is an integral operator from $L^2(\pi)$ to $L^2(\pi)$ such that

$$(Kf)(s) = \int k(s, t) f(t) \pi(t) dt, \quad (4)$$

with kernel

$$k(s, t) = E[h_j(s; \theta_0) \overline{h_j(t; \theta_0)}] = \psi_0(s - t; \theta_0) - \psi_0(s; \theta_0) \psi_0(-t; \theta_0), \quad (5)$$

and its relation operator R , which is an integral operator from $L^2(\pi)$ to $L^2(\pi)$ with kernel

$$r(s, t) = E[h_j(s; \theta_0) h_j(t; \theta_0)] = k(s, -t). \quad (6)$$

In the sequel, we denote by λ_j and ϕ_j the eigenvalues and orthonormal eigenfunctions of K , respectively, which are solutions to the functional equation $(K\phi_j)(t) = \lambda_j\phi_j(t)$. The ϕ_j are not uniquely defined because one can multiply a complex function by a complex number on the unit circle without altering its norm. However, our test statistic is invariant to this. Figures B2a and B2c in the Supplemental Appendix present the eigenfunctions associated with the largest two eigenvalues for the covariance operator K for the standard normal when the weighting function π is itself a normal with zero mean and scale parameter ω for two values of ω . In turn, Figures B2b and B2d show the corresponding operator of the standardized uniform distribution on $(-\sqrt{3}, \sqrt{3})$ for the same Gaussian weighting function. As can be seen in these figures, if we arrange the eigenvalues in decreasing order, the eigenfunctions associated with even (odd) eigenvalues are even (odd) functions in these two examples. We also report in Figures B2e and B2f the largest five eigenvalues for those distributions. As we shall see below, the main effect of changing ω will be to change the relative weights given to small and large values of the CF argument t .

We are interested in applying Hansen (1982)'s J test of overidentifying restrictions to our continuum of moments. To illustrate the difficulties that may arise, assume for a moment that $\hat{h}(\theta)$ is a finite dimensional m -vector obtained from a rough discretization of \mathbb{R}^q , so that $\sqrt{n}\hat{h}(\theta_0) \xrightarrow{d} \mathcal{N}(0, \mathcal{K})$ and \mathcal{K} is a nonsingular $m \times m$ matrix. Assuming for simplicity that both \mathcal{K} and θ are known, the usual J test for overidentifying restrictions is

$$J = n\hat{h}^*(\theta)\mathcal{K}^{-1}\hat{h}(\theta), \quad (7)$$

where $*$ denotes the complex conjugate transpose of a vector/matrix. Now if we let m grow by taking a denser and denser grid, then the matrix \mathcal{K} becomes increasingly ill-conditioned, in the sense that the ratio of its largest eigenvalue to its smallest one increases dramatically, so \mathcal{K}^{-1} may be numerically unreliable for large m .

In our setting, the covariance matrix \mathcal{K} is replaced by the aforementioned covariance operator K (see Supplemental Appendix A.1), which has a countable infinite number of positive eigenvalues λ_j , $j = 1, 2, \dots$ (arranged in decreasing order) and associated eigenfunctions ϕ_j . As we will see later, this operator is compact, meaning that its inverse is not bounded. Consequently, its smallest eigenvalues will converge to zero as j goes to infinity, so taking the inverse of K is problematic. In terms of the spectral decomposition of K , the direct analogue to the J test statistic in (7) would be written as

$$\left\langle \sqrt{n}\hat{h}, K^{-1}\sqrt{n}\hat{h} \right\rangle = \sum_j \frac{1}{\lambda_j} \left| \left\langle \sqrt{n}\hat{h}, \phi_j \right\rangle \right|^2 \quad (8)$$

where the dependence on θ is omitted for simplicity and $|\cdot|$ denotes the modulus of complex numbers. This expression will blow up because of the division by the small eigenvalues λ_j for large j . This is related to the problem of solving an integral equation $Kf = g$ where g is known and f is the object of interest. This problem is said to be ill-posed because f is not continuous in g . Indeed, a small perturbation in g will result in a large change in f . To stabilize the solution, one needs to use some regularization scheme (see Kress (1999) and Carrasco, Florens, and Renault (2007) for various possibilities). As in Carrasco and Florens (2000), we use Tikhonov regularization, which consists in replacing $K^{-1}g$ by the regularized solution $(K^2 + \alpha I)^{-1}Kg$ where $\alpha \geq 0$ is a regularization parameter. We use the notation $(K^\alpha)^{-1}$ for $(K^2 + \alpha I)^{-1}K$, which is the operator with eigenvalues $\frac{\lambda_j}{\lambda_j^2 + \alpha}$ and corresponding eigenfunctions ϕ_j , and $(K^\alpha)^{-1/2}$ for the operator with eigenvalues $\frac{\sqrt{\lambda_j}}{\sqrt{\lambda_j^2 + \alpha}}$ and the same eigenfunctions.

Thus, the regularized version of the J test is

$$\left\| (K^\alpha)^{-1/2} \sqrt{n} \hat{h} \right\|^2 = \sum_j \frac{\lambda_j}{\lambda_j^2 + \alpha} \left| \left\langle \sqrt{n} \hat{h}, \phi_j \right\rangle \right|^2. \quad (9)$$

Comparing the expressions (8) and (9), we observe that $\frac{1}{\lambda_j}$ has been replaced by $\frac{\lambda_j}{\lambda_j^2 + \alpha}$, which is bounded.

We will consider various versions of this test depending on whether:

- θ_0 is known or estimated,
- K is known or estimated,
- α is fixed or goes to zero.

Consider the case where α is fixed; if we are willing to assume that θ_0 is known, so that the distribution under the null hypothesis is completely specified and the operator K is known, then the first test we should consider is

$$J(\theta_0, K) = \sum_j \frac{\lambda_j}{\lambda_j^2 + \alpha} \left| \left\langle \sqrt{n} \hat{h}, \phi_j \right\rangle \right|^2. \quad (10)$$

As we explain in Appendix A.1, the test statistic (10) can be arbitrarily well approximated from a numerical point of view by a regularized version of the matrix expression (7). Specifically, if we evaluated the CF at a very fine but discrete grid of m points over a finite range of values of the argument t , then

$$J(\theta_0, K) = n \hat{\underline{h}}(\theta_0)^* \left(\frac{\mathcal{K}}{m} \right)^{1/2} \left[\left(\frac{\mathcal{K}}{m} \right)^2 + \alpha I \right]^{-1} \left(\frac{\mathcal{K}}{m} \right)^{1/2} \hat{\underline{h}}(\theta_0). \quad (11)$$

Several issues related to the practical implementation of this test (in particular the computation of the eigenvalues of K) are discussed in the Supplemental Appendix A.1.

When θ is unknown, however, the operator K is only known up to θ . Let $\tilde{\theta}$ be a consistent estimator of θ obtained for instance from

$$\tilde{\theta} = \arg \min_{\theta \in \Theta} \left\| \hat{h}(\cdot; \theta) \right\|^2.$$

In this context, the integral operator $K_{\tilde{\theta}}$ can be defined as in (4) but with kernel

$$k(s, t) = \psi_0(s - t; \tilde{\theta}) - \psi_0(s; \tilde{\theta})\psi_0(-t; \tilde{\theta}).$$

Let $\{\lambda_{j\tilde{\theta}}, \phi_{j\tilde{\theta}}\}$ $j = 1, \dots, m$ be the eigenvalues and eigenfunctions of the operator $K_{\tilde{\theta}}$. Then the second test we consider is

$$J(\hat{\theta}, K_{\tilde{\theta}}) = \sum_j \frac{\lambda_{j\tilde{\theta}}}{\lambda_{j\tilde{\theta}}^2 + \alpha} \left| \left\langle \sqrt{n}\hat{h}(\cdot; \hat{\theta}), \phi_{j\tilde{\theta}} \right\rangle \right|^2 = \min_{\theta \in \Theta} \sum_j \frac{\lambda_{j\tilde{\theta}}}{\lambda_{j\tilde{\theta}}^2 + \alpha} \left| \left\langle \sqrt{n}\hat{h}(\cdot; \theta), \phi_{j\tilde{\theta}} \right\rangle \right|^2$$

where $\hat{\theta}$ corresponds to the argument of the minimization.

Alternatively, we may prefer to estimate K using a sample covariance operator. In fact, there are two obvious possibilities. The first one is to use the integral estimator \hat{K} with uncentered kernel

$$\hat{k}(s, t) = \frac{1}{n} \sum_{i=1}^n h_i(s; \tilde{\theta})h_i(-t; \tilde{\theta}),$$

where $\tilde{\theta}$ is a consistent first step estimator of θ . On the other hand, the second possibility is the integral operator \hat{K} with centered kernel

$$\hat{k}(s, t) = \frac{1}{n} \sum_{i=1}^n h_i(s)h_i(-t),$$

where

$$h_i(s) = h_i(s; \theta) - \hat{h}(s; \theta) = e^{isX_i} - \frac{1}{n} \sum_{l=1}^n e^{isX_l}.$$

The advantage of the second estimator is that it does not require a first step estimator of θ and thereby it may be more robust to misspecification.

For computational reasons, it is convenient to rewrite the test statistics (9), which use as eigenvalues and eigenfunctions those of \hat{K} and \hat{K} , in terms of certain matrices and vectors (see Carrasco et al (2007) for analogous expressions for \hat{K} under time series dependence). Specifically, we obtain the following two expressions:

i) The test based on \hat{K} , which can be computed as

$$J(\hat{\theta}, \hat{K}_{\hat{\theta}}) = \min_{\theta \in \Theta} \underline{v}(\theta)^* [\alpha I + C^2]^{-1} \underline{v}(\theta) \quad (12)$$

where $\underline{v}(\theta)$ is a $n \times 1$ vector with l -th element $v_l(\theta) = \int \overline{h_l(t; \tilde{\theta})} \hat{h}(t; \theta) \pi(t) dt$, C is an $n \times n$ matrix with (i, l) element c_{il}/n with $c_{il} = \langle h_l(t; \tilde{\theta}), h_i(t; \tilde{\theta}) \rangle$ (see Supplemental Appendix A.2 for analytical expressions for these integrals).

ii) The test based on \hat{K} , whose matrix expression is

$$J(\hat{\theta}, \hat{K}) = \min_{\theta \in \Theta} \hat{\underline{v}}(\theta)^* [\alpha I + \hat{C}^2]^{-1} \hat{\underline{v}}(\theta) \quad (13)$$

where $\hat{\underline{v}}(\theta)$ is a $n \times 1$ vector with l -th element $\hat{v}_l(\theta) = \int \overline{h_l(t)} \hat{h}(t; \theta) \pi(t) dt$, \hat{C} is an $n \times n$ matrix with (i, l) element \hat{c}_{il}/n with $\hat{c}_{il} = \langle h_l(t), h_i(t) \rangle$, and $\hat{\theta}$ is the argument of the minimization. Note that $\hat{C} = (I - \ell\ell'/n)C(I - \ell\ell'/n)$, where ℓ is a vector of n ones.

In Sections 3 and 4, we will study the asymptotic distribution of the test statistics $J(\theta_0, K)$, $J(\hat{\theta}, K_{\hat{\theta}})$, $J(\hat{\theta}, \hat{K}_{\hat{\theta}})$ and $J(\hat{\theta}, \hat{K})$ and show that they converge under H_0 to a weighted sum of χ^2 's whose weights depend on θ . Given the eigenvalues, those weights and hence their asymptotic distributions are known, so we can compute the p-value of these quadratic forms in normal variables using the approach in Imhof (1961). Nevertheless, we rely on the parametric bootstrap in the simulations to improve the small sample properties of our proposed procedures.

For all the tests presented so far, α is fixed, so that our regularized inverse $(K^\alpha)^{-1}$ is a biased approximation of K^{-1} . It is possible to approach K^{-1} by letting α go to zero at a suitable rate. However, a test based on (9) with α going to zero would tend to a chi-square with infinite degrees of freedom, and hence diverge. For that reason, we explain next how to center and rescale it following Carrasco and Florens (2000). Let $h_j(t; \theta_0)$ denote the influence function (3) evaluated at the true θ_0 (here θ_0 is assumed to be known to simplify). Similarly, let $\hat{\lambda}_j$ denote the eigenvalues of \hat{K} , the sample covariance operator of $h_j(t; \theta_0)$,

$$\hat{a}_j = \frac{\hat{\lambda}_j^2}{\hat{\lambda}_j^2 + \alpha}, \quad \hat{p}_n = \sum_{j=1}^n \hat{a}_j, \quad \text{and} \quad \hat{q}_n = 2 \sum_{j=1}^n \hat{a}_j^2. \quad (14)$$

After appropriate centering and rescaling, we obtain:

$$J_{\alpha_n} = \frac{\left\| (\hat{K}^{\alpha_n})^{-1/2} \sqrt{n} \hat{h}(\cdot; \theta_0) \right\|^2 - \hat{p}_n}{\sqrt{\hat{q}_n}}. \quad (15)$$

In Section 5, we show that J_{α_n} converges to a standard normal distribution under the null.

3 J test when α is fixed and the parameter is known

3.1 Distribution under local alternatives

The $J(\theta_0, K)$ statistic in (10) with α fixed is part of a larger class of tests based on weighted L^2 statistics that we will denote by T_B in the sequel. Let B be a nonrandom bounded linear operator from $L^2(\pi)$ to $L^2(\pi)$ and B_n a sequence of random bounded linear operators from $L^2(\pi)$ to $L^2(\pi)$ such that $\|B_n - B\| \xrightarrow{P} 0$ as n goes to infinity, where $\|\cdot\|$ is the sup-norm. Assume moreover that the null space of B equals $\{0\}$; otherwise the test would lack power against certain alternatives. Popular choices of B satisfying our assumptions include $B = I$ as in Epps and Pulley (1983), Bierens and Wang (2012) and Leucht (2012), as well as $B = (K^\alpha)^{-1/2}$ with $\alpha > 0$ fixed. Note that B is not necessarily real.

In this section and the next one we focus on tests based on weighted L^2 statistics

$$T_B = \left\| B_n \sqrt{n} \hat{h} \right\|^2 = \int \left| B_n \sqrt{n} \hat{h} \right|^2(t) \pi(t) dt, \quad (16)$$

where $\hat{h}(t) = \sum_{j=1}^n [e^{itX_j} - \psi_0(t)]$ and $\psi_0(t) = \psi_0(t; \theta_0)$.

First we express the null and alternative hypotheses in terms of the density function. Specifically, let f_0 be a density with respect to Lebesgue measure (the extension to the case of another measure for discrete or mixed random variables is straightforward and will not be treated to avoid cumbersome notation), then:

$$H_0 : f(x) = f_0(x),$$

$$H_{1n}(c) : f_n(x) = f_0(x) \left[1 + \frac{cu(x)}{\sqrt{n}} \right]$$

where c is a scalar and u is such that $\int u(x) f_0(x) dx = 0$ so that f_n integrates to 1.

Given that there is a one-to-one mapping between the density and the CF through the Fourier inversion theorem (see (1)), we can reformulate H_0 and H_{1n} in terms of the CF instead. Thus, we obtain

$$H_0 : \psi = \psi_0,$$

$$H_{1n}(c) : \psi_n = \psi_0 + \frac{c\eta}{\sqrt{n}}$$

where $\eta(t) = \int e^{itx} u(x) f_0(x) dx$.

To guarantee the uniqueness of the representation, η needs to be normalized. Many normalizations could be used. For convenience, we impose the normalization condition $\|u\|_{L^2(f_0)} = E[u^2(X)] = 1$. Remark that by construction, $\eta(0) = 0$ and $\overline{\eta(t)} = \eta(-t)$. Moreover given $|\psi_n| \leq 1$, η is bounded. In this context, η represents the direction of the alternative, while c represents the distance from the null.

First, we establish some results on the operator K of form (4) with kernel (5), suppressing

the dependence on θ_0 for simplicity.

Lemma 1 *K is a self-adjoint positive definite Hilbert-Schmidt operator from $L^2(\pi)$ to $L^2(\pi)$ and the sum of its eigenvalues is bounded by 1.*

Lemma 1 implies two things: that K has a countable spectrum and that the sum of its eigenvalues is less than 1.

Example *Consider the CF of a univariate normal with mean μ and variance σ^2 ; it turns out that when using a normal weighting function with zero mean and scale parameter ω , we can obtain analytical solutions for the sums of both λ 's and λ^2 's. Specifically, the expressions are*

$$\sum_j \lambda_j = 1 - \frac{1}{\sqrt{1 + 2\sigma^2\omega^2}}$$

and

$$\sum_j \lambda_j^2 = \frac{1}{\omega + 2\sigma^2\omega^3} \left(\frac{e^{-\frac{4\mu^2\omega^2}{1+2\sigma^2\omega^2}}}{\sqrt{1 + 2\sigma^2\omega^2}} + \frac{e^{-\frac{4\mu^2\omega^2}{1+4\sigma^2\omega^2}}}{\sqrt{1 + 4\sigma^2\omega^2}} - \frac{2}{\sqrt{1 + 4\sigma^2\omega^2 + 3\sigma^4\omega^4}} \right).$$

As can be seen from the above expressions, the sums of both λ 's and λ^2 's depend on the scale ω of the weighting function.

Assumption 1 $X_i, i = 1, 2, \dots$ are iid.

Proposition 2 *Assume that Assumption 1 holds. Under H_{1n} , as n goes to infinity*

$$\sqrt{n\hat{h}} \Rightarrow \mathcal{CN}(c\eta, K, R)$$

in $L^2(\pi)$ where the covariance operator K is the integral operator with kernel (5) and the relation operator R is the integral operator with kernel (6).

To establish the asymptotic distribution of $\|B_n\sqrt{n\hat{h}}\|^2$, it is useful to relate the complex process $B_n\sqrt{n\hat{h}}$ to a real process following an approach similar to that of Ducharme, Lafaye de Micheaux and Marchina (2016). Let

$$Z_n = \begin{pmatrix} B_n\sqrt{n\hat{h}} \\ \overline{B_n\sqrt{n\hat{h}}} \end{pmatrix}, Y_n = \begin{pmatrix} \operatorname{Re} B_n\sqrt{n\hat{h}} \\ \operatorname{Im} B_n\sqrt{n\hat{h}} \end{pmatrix}, \text{ and } M = \frac{1}{2} \begin{pmatrix} 1 & 1 \\ -i & i \end{pmatrix}.$$

Before stating the general results, note that Z_n and Y_n are elements of $(L^2(\pi))^2$, the space of 2×1 vectors of complex valued functions with inner product denoted $\langle f, g \rangle_2$ and defined as:

$$\langle f, g \rangle_2 = \int_{-\infty}^{\infty} f_1(t) \overline{g_1(t)} \pi(t) dt + \int_{-\infty}^{\infty} f_2(t) \overline{g_2(t)} \pi(t) dt$$

for $f = (f_1, f_2)'$ and $g = (g_1, g_2)'$. On this space, the norm is denoted as $\|f\|_2 = \sqrt{\langle f, f \rangle_2}$.

Note that \widehat{h} is not an arbitrary element of $L^2(\pi)$. By definition, \widehat{h} involves a Fourier transform of a real valued process and hence satisfies the property $\overline{\widehat{h}(t)} = \widehat{h}(-t)$. Its covariances satisfy $E\left(h_j(s)\overline{h_j(t)}\right) = \psi_0(s-t) - \psi_0(s)\psi_0(-t) \equiv k(s,t)$, $E(h_j(s)h_j(t)) = \psi_0(s+t) - \psi_0(s)\psi_0(t) = k(s,-t) \equiv r(s,t)$, $E\left(\overline{h_j(s)}h_j(t)\right) = \overline{k(s,t)} = k(-s,-t)$, and $E\left(\overline{h_j(s)}\overline{h_j(t)}\right) = \overline{r(s,t)} = k(s,-t)$. There is one-to-one mapping between the complex process Z_n and the real process Y_n through $Z_n = M^{-1}Y_n$. Moreover, $\|B_n\sqrt{n}\widehat{h}\|_2^2 = \|Y_n\|_2^2$. As Y_n converges to a bivariate real Gaussian process, then $\|Y_n\|_2^2$ converges to a weighted sum of chi-square distributions.

We have the relationship

$$M^{-1} = \begin{pmatrix} 1 & i \\ 1 & -i \end{pmatrix} = 2M^*$$

where M^* is the adjoint of M on $(L^2(\pi))^2$. Let Γ be the covariance operator of $\begin{pmatrix} B\sqrt{n}\widehat{h} \\ \overline{B}\sqrt{n}\overline{\widehat{h}} \end{pmatrix}$ under H_0 . It is an integral operator from $(L^2(\pi))^2$ to $(L^2(\pi))^2$ such that

$$g = \begin{pmatrix} g_1 \\ g_2 \end{pmatrix} \in (L^2(\pi))^2 \rightarrow \Gamma g = \begin{pmatrix} BKB^* & BR\overline{B}^* \\ \overline{B}\overline{R}B^* & \overline{B}\overline{K}\overline{B}^* \end{pmatrix} g$$

where \overline{K} and \overline{R} are the integral operators with kernels $\overline{k}(s,t) = k(-s,-t)$ and $\overline{r}(s,t) = k(s,-t)$, respectively. The covariance operator of Y_n is $M\Gamma M^*$.

Let b_j, ζ_j , for $j = 1, 2, \dots$, be the nonzero eigenvalues (arranged in decreasing order and repeated according to their order of multiplicity) and 2×1 eigenfunctions of $M\Gamma M^*$. Let ζ_l be the eigenfunctions of $M\Gamma M^*$ associated with the 0 eigenvalue of $M\Gamma M^*$ when Γ is singular. As $M\Gamma M^*$ is real, so are b_j, ζ_j , and ζ_l . Lemma 13 in Appendix gives a complete characterization of ζ_j and ζ_l . In particular, it shows that when $B = I$, then one can associate to each eigenfunction ϕ_j of K (corresponding to the eigenvalue λ_j), a pair of 2×1 functions $\zeta_j(t) = M(\phi_j(t), \phi_j(-t))'$ and $\zeta_l(t) = M(\phi_l(t), -\phi_l(-t))'$ (with $l = j$) which are the eigenfunctions of $M\Gamma M^*$ associated with the eigenvalues λ_j and 0 respectively. As these eigenfunctions are orthogonal, they form a complete orthonormal basis of $(L^2(\pi))^2$. Then, the asymptotic distribution of T_B follows from Karhunen-Loeve theorem (see Lemmas 12-14 in Appendix for more details). Further, let $\tilde{\eta} = (\text{Re}(B\eta), \text{Im}(B\eta))'$.

Proposition 3 *Assume that Assumption 1 holds.*

(a) *Under H_{1n} , we have*

$$T_B \xrightarrow{d} \sum_{j=1}^{\infty} b_j \chi_j^2(1, \delta_j) + c \sum_l \langle \tilde{\eta}, \zeta_l \rangle_2^2 = \sum_{j=1}^{\infty} b_j \left(e_j + \frac{c \langle \tilde{\eta}, \zeta_j \rangle_2}{\sqrt{b_j}} \right)^2 + c \sum_l \langle \tilde{\eta}, \zeta_l \rangle_2^2$$

where $\chi_j^2(1, \delta_j)$, $j = 1, 2, \dots$ denote independent noncentral chi-square random variables with 1 degree of freedom and non centrality parameter $\delta_j = c^2 \langle \tilde{\eta}, \zeta_j \rangle^2 / b_j$ while e_j , $j = 1, 2, \dots$ are the underlying independent standard normal variables.

(b) If moreover π is symmetric around 0 and B is either the identity operator or $B = (K^\alpha)^{-1/2}$, a simplification yields the following distribution under H_{1n} :

$$T_B \stackrel{d}{\rightarrow} \sum_{j=1}^{\infty} a_j \chi_j^2(1, \delta_j) = \sum_{j=1}^{\infty} a_j \left(e_j + \frac{c \langle B\eta, \phi_j \rangle}{\sqrt{a_j}} \right)^2$$

where $\chi_j^2(1, \delta_j)$, $j = 1, 2, \dots$ denote independent noncentral chi-square random variables with 1 degree of freedom and non centrality parameter $\delta_j = c^2 \langle B\eta, \phi_j \rangle^2 / a_j$ where a_j , ϕ_j , $j = 1, 2, \dots$ are the eigenvalues (arranged in decreasing order) and eigenfunctions of BKB^* , while e_j , $j = 1, 2, \dots$ are iid $\mathcal{N}(0, 1)$.

Remark 1 While part (a) of Proposition 3 holds for general moment conditions and arbitrary π , the proof of part (b) heavily relies on the fact that the complex conjugate of $\hat{h}(t)$ equals $\hat{h}(-t)$ and the symmetry of π . We observe that, in this case, the asymptotic distribution under H_0 is a weighted sum of chi-squares as we would obtain if the moment conditions were real.

Remark 2 The previous proposition will not apply if $B = K^{-1/2}$. In that case, B is not bounded, which violates one of the assumptions. Moreover, $\mathcal{N}(0, I)$ is not a Gaussian process because the trace of its covariance operator (the identity operator) is infinite. We will discuss the case $B = (K^\alpha)^{-1/2}$ when α goes to zero in Section 5.

Remark 3 Now we comment on the cases $B = I$ and $B = (K^\alpha)^{-1/2}$. Recall that the case $B = I$ corresponds to Bierens and Wang (2012) test whereas $B = (K^\alpha)^{-1/2}$ corresponds to a J-test where the weighting operator is a biased estimator of K^{-1} . We see that as soon as $\langle B\eta, \phi_j \rangle \neq 0$ for some j , the test statistic T_B will have non trivial power. But because $\{\phi_j\}$ forms an orthonormal basis of $L^2(\pi)$, then $B\eta = \sum_j \langle B\eta, \phi_j \rangle \phi_j$ and by Parseval's identity, $\|B\eta\|^2 = \sum_j \langle B\eta, \phi_j \rangle^2 > 0$. It follows that $\langle B\eta, \phi_j \rangle$, $j = 1, 2, \dots$ cannot all be zero simultaneously. Therefore, T_B has indeed non trivial power against all local alternatives of the form H_{1n} , and against all fixed alternatives a fortiori. However, if $\langle B\eta, \phi_j \rangle^2$ is small (as will be the case for most j since the sequence $\langle B\eta, \phi_j \rangle^2$ is summable), the power against local alternatives in the j th direction may be poor.

Remark 4 Now, we consider the power of T_B against a fixed alternative of the form $H_1 : \psi_n = \psi_0 + \eta$. We have

$$\frac{T_B}{n} = \int (B_n h_n)^2(t) \pi(t) dt \xrightarrow{P} \|B\eta\|^2.$$

Hence,

$$\begin{aligned}\frac{T_I}{n} &\xrightarrow{P} \|\eta\|^2 = \sum_j \langle \eta, \phi_j \rangle^2, \\ \frac{T_{(K_n^\alpha)^{-1/2}}}{n} &\xrightarrow{P} \left\| (K^\alpha)^{-1/2} \eta \right\|^2 = \sum_j \frac{\lambda_j}{\lambda_j^2 + \alpha} \langle \eta, \phi_j \rangle^2.\end{aligned}$$

As $\lambda_j < 1$, we can always find an α sufficiently small so that $\frac{\lambda_j}{\lambda_j^2 + \alpha} > 1$ for all j , and hence $T_{(K_n^\alpha)^{-1/2}}$ has better power than T_I in the sense that $p \lim \frac{T_{(K_n^\alpha)^{-1/2}}}{n} > p \lim \frac{T_I}{n}$ against any fixed alternative.¹

In the next subsection, we will study the power properties of these tests in more detail.

3.2 Alternatives with maximum power

It is well-known that there is no uniformly most powerful test for assessing H_0 and that goodness-of-fit tests have good power only against certain local alternatives (see Neuhaus (1976), Janssen (2000), Escanciano (2009), and Lehmann and Romano (2005, Section 14.6)). In this subsection, we will characterize the alternative with maximum power.

Let $L^2(f_0) < \infty$ denote the L^2 space of real functions $\varphi(X)$ such that we can define $\|\varphi\|_{L^2(f_0)}^2 = \int \varphi^2(x) f_0(x) dx$.

In this section, we focus on the case where $B = I$ or $(K^\alpha)^{-1/2}$, so we will assume that the following assumption holds.

Assumption 1' $X_i, i = 1, 2, \dots$ are iid, π is symmetric around 0 and $B = I$ or $(K^\alpha)^{-1/2}$.

In this set up, we define the asymptotic local power function $\Pi_B(a, c, u)$ as

$$\Pi_B(a, c, u) = \lim_{n \rightarrow \infty} P[T_B \geq c_a | H_{1n}(c)],$$

where c_a is the critical value such that T_B achieves a level a , i.e. $\lim_{n \rightarrow \infty} P(T_B \geq c_a | H_0) = a$.

To analyze the power of these statistics, it is useful to rewrite K as T^*T , where T is an operator from $L^2(\pi)$ to $L^2(f_0)$ and T^* is the adjoint operator from $L^2(f_0)$ to $L^2(\pi)$. Such a decomposition has been used to study the power of Cramer von Mises type tests by Neuhaus (1976, equation (1.9)) and Escanciano (2009, p.168).

The operators T and T^* are as follows:

$$\begin{aligned}T &: L^2(\pi) \rightarrow L^2(f_0), \\ (T\phi)(X) &= \int \overline{h(X; t)} \phi(t) \pi(t) dt, \\ T^* &: L^2(f_0) \rightarrow L^2(\pi), \text{ and}\end{aligned}$$

¹We thank a referee for suggesting this remark.

$$(T^*\varphi)(t) = \int h(x;t) \varphi(x) f_0(x) dx.$$

Moreover, TB^* is compact and admits a singular system $\{\sqrt{a_j}, \phi_j, \varphi_j\}$, where $TB^*\phi_j = \sqrt{a_j}\varphi_j$ and $BT^*\varphi_j = \sqrt{a_j}\phi_j$. Therefore, ϕ_j are the eigenfunctions of $BT^*TB^* = BKB^*$ and φ_j those of B^*TT^*B . Thus, φ_j can be interpreted as principal components.

Observe that $\eta = T^*u$. Indeed, if we use the property of Fourier transforms and the fact that $\int u(x) f_0(x) dx = 0$, we will have that

$$\begin{aligned} (T^*u)(t) &= \int [e^{itx} - \psi_0(t)]u(x) f_0(x) dx \\ &= \int e^{itx}u(x) f_0(x) dx - \psi_0(t) \int u(x) f_0(x) dx = \eta(t). \end{aligned}$$

Hence, the relation $\eta = T^*u$ implies that

$$\frac{\langle B\eta, \phi_j \rangle}{\sqrt{a_j}} = \frac{\langle BT^*u, \phi_j \rangle}{\sqrt{a_j}} = \frac{\langle u, TB^*\phi_j \rangle}{\sqrt{a_j}} = \langle u, \varphi_j \rangle. \quad (17)$$

From (17) and Proposition 2, it follows that under $H_{1n}(c)$,

$$T_B \xrightarrow{d} \sum_{j=1}^{\infty} a_j (e_j + c \langle u, \varphi_j \rangle)^2. \quad (18)$$

Note that the sequence $\{\varphi_j\}$, for $j = 1, 2, \dots$, forms a complete orthonormal basis of $\overline{\mathcal{R}(TB^*)} = L^2(f_0) \cap \{u : E(u) = 0\}$. Hence, the alternatives of interest are linear combinations of the eigenfunctions φ_j 's. In this context, the analysis of the limiting distribution in (18) and the orthogonality of the φ_j 's allow us to establish the following result:

Proposition 4 *Assume that Assumption 1' holds. The limiting power $\Pi_B(\alpha, c, u)$ has the following properties.*

- (a) $\max_u \{\Pi_B(a, c, u) : u \in L^2(f_0), E(u) = 0, \|u\|_{L^2(f_0)} = 1\} = \Pi_B(a, c, \varphi_1)$,
- (b) $\Pi_B(a, c, \varphi_j) \leq \Pi_B(a, c, \varphi_i)$ for $1 \leq i \leq j$,
- (c) $\lim_{j \rightarrow \infty} \Pi_B(a, c, \varphi_j) = a$.

Proposition 4 says that (a) the maximum power is achieved for the local alternative $u = \varphi_1$ corresponding to the first principal component, (b) the power decreases when considering higher-order principal components, (c) finally, the power goes down to the level of the test, a , for the highest frequency (case $j \rightarrow \infty$).

As we saw before, in general φ_j depends on B , so that the alternative with maximum power will be different for different tests T_B .

But if we consider more specifically the cases $B = I$ and $B = (K^\alpha)^{-1/2}$, the φ_j are the same because they correspond to the eigenfunctions of TT^* . Hence, the alternative for which the tests

T_B for $B = I$ and $B = (K^\alpha)^{-1/2}$ are the most powerful coincides, and corresponds to $\eta = \varphi_1$.

When $B = I$, then $a_j = \lambda_j$, i.e. the eigenvalues of K decline quickly towards 0. So the test T_B with $B = I$ concentrates its power on the first principal component. On the other hand, when $B = (K^\alpha)^{-1/2}$, $a_j = \frac{\lambda_j^2}{\lambda_j^2 + \alpha}$ instead will decline slower towards 0 for smaller α . Consequently, power will be more evenly spread among the first few directions when $B = (K^\alpha)^{-1/2}$ than when $B = I$. This point is illustrated graphically in the case of the normal distribution. Figure 1 shows the decline of λ_j and a_j while Figure 2 reports the first three alternatives with maximum power. In turn, Figure 3 plots the asymptotic powers of T_I and $T_{(K_n^\alpha)^{-1/2}}$ for different values of α . In the extreme case where $\alpha = 0$, we would have $a_j = 1$, which means that power would be evenly spread among all alternatives. However, in this case the null distribution is a Chi-square with infinite degrees of freedom and the resulting test has power equal to size for any local alternative; see Lemma 14.3.1 of Lehmann and Romano (2005). We will consider the case where $\alpha \rightarrow 0$ in greater detail in Section 5.

4 J test when α is fixed and the parameter is unknown

4.1 Distribution under local alternatives

Let $\Theta \subset \mathbb{R}^p$ be the parameter space of θ . Let $\theta_0 \in \Theta$ be the true value of θ under H_0 . Consider

$$H_0 : f(x) = f_0(x; \theta_0),$$

$$H_{1n}(c) : f_n(x) = f_0(x; \theta_0) \left[1 + \frac{cu(x)}{\sqrt{n}} \right]$$

where c is a scalar and u is such that $\int u(x) f_0(x) dx = 0$. Equivalently, the hypotheses can be written as

$$H_0 : \psi = \psi_0(\cdot; \theta_0),$$

$$H_{1n}(c) : \psi_n = \psi_0(\cdot; \theta_0) + \frac{c\eta}{\sqrt{n}}$$

where $\eta(t) = \int e^{itx} u(x) f_0(x; \theta_0) dx$.

Assumption 2 Under H_{1n} , $\|B_n - B\| \xrightarrow{P} 0$. Under H_1 , $\|B_n - B_1\| \xrightarrow{P} 0$ where both B and B_1 are bounded linear operators and B_1 may differ from B . The null spaces of B and B_1 equal $\{0\}$.

In the sequel, we denote by P_0 the law of X_i under H_0 , P_n the law of X_i under H_{1n} , and P_1 the law of X_i under H_1 .

Assumption 3 P_n is contiguous to P_0 .

This condition is standard in the goodness-of-fit literature and imposes some mild restrictions on the density. Sufficient conditions for this assumption to be true are given in Lehmann and Romano (2005). They also provide a variety of examples.

Assumption 4 The parameter space Θ is a compact subset of \mathbb{R}^p . The true parameter θ_0 is contained in the interior of Θ . $\psi_0(\tau; \theta)$ is continuously differentiable with respect to θ .

Assumption 5 (identification) $\psi_0(\tau; \theta) = \psi_0(\tau; \theta_0)$ for all $\tau \Leftrightarrow \theta = \theta_0$.

Let

$$\hat{\theta} = \arg \min_{\theta \in \Theta} \left\| B_n \hat{h}(\cdot; \theta) \right\|$$

and define

$$D_0 = \left. \frac{\partial \psi_0(\cdot; \theta)}{\partial \theta} \right|_{\theta = \theta_0}.$$

Note that the result in Proposition 2 remains valid here. Namely, $\sqrt{n} \hat{h}(\theta_0) \Rightarrow \mathcal{N}(c\eta, K)$ under H_{1n} , where K is an integral operator with kernel $k(s, t) = \psi(s - t; \theta_0) - \psi(s; \theta_0) \psi(-t; \theta_0)$.

Proposition 5 Suppose Assumptions 1-5 hold. Under H_0 , $\hat{\theta}$ is a consistent estimator of θ_0 and

$$\sqrt{n}(\hat{\theta} - \theta_0) \xrightarrow{d} \mathcal{N}(0, \langle BD_0, BD_0 \rangle^{-1} \langle BD_0, (BKB^*) BD_0 \rangle \langle BD_0, BD_0 \rangle^{-1}).$$

Moreover, under H_1 ,

$$\hat{\theta} \xrightarrow{P_1} \theta_1 = \arg \min_{\theta \in \Theta} \left\| B_1 E^{P_1} (h_j(\cdot; \theta)) \right\|.$$

Let L be the operator from $L^2(\pi)$ to $L^2(\pi)$ such that for all $\varphi \in L^2(\pi)$

$$(L\varphi)(\tau) = \varphi(\tau) - D_0(\tau) \langle BD_0, BD_0 \rangle^{-1} \langle B^* BD_0, \varphi \rangle.$$

Let \tilde{K} be the integral operator from $L^2(\pi)$ to $L^2(\pi)$ with kernel

$$\begin{aligned} \tilde{k}(s, t) &= k(s, t) - D_0(s) \langle BD_0, BD_0 \rangle^{-1} (KB^* BD_0)(t) \\ &\quad - \overline{D_0(t)} \langle BD_0, BD_0 \rangle^{-1} \overline{(KB^* BD_0)(s)} \\ &\quad + D_0(s) \langle BD_0, BD_0 \rangle^{-1} \langle BD_0, (BKB^*) BD_0 \rangle \langle BD_0, BD_0 \rangle^{-1} \overline{D_0(t)}' \end{aligned}$$

In addition, let $\tilde{a}_j, \tilde{\phi}_j$, for $j = 1, 2, \dots, J$, denote the eigenvalues (arranged in decreasing order) and eigenvectors of $B\tilde{K}B^*$. Finally, define $\tilde{\delta}_j = c^2 \langle BL\eta, \tilde{\phi}_j \rangle^2 / \tilde{a}_j$.

Proposition 6 Suppose Assumptions 1', 2-5 hold. Under H_{1n} , we have

(i) $\sqrt{n} \hat{h}(\hat{\theta})$ converges to a complex Gaussian process with mean $cL\eta$ and covariance operator \tilde{K} in $L^2(\pi)$.

(ii)

$$T_B \xrightarrow{d} \sum_{j=1}^{\infty} \tilde{a}_j \chi_j^2(1, \tilde{\delta}_j) = \sum_{j=1}^{\infty} \tilde{a}_j \left(e_j + \frac{c \langle BL\eta, \tilde{\phi}_j \rangle}{\sqrt{\tilde{a}_j}} \right)^2$$

where $\chi_j^2(1, \tilde{\delta}_j)$, $j = 1, 2, \dots$ denote independent noncentral chi-square r.v. with 1 degree of freedom and non centrality parameter $\tilde{\delta}_j$ and e_j , $j = 1, 2, \dots$ are the underlying independent standard normal variables.

Proposition 6 implies that T_B has non trivial power against all local alternatives η for which $L\eta \neq 0$, i.e. those η such that $\eta \neq v'D_0$, where v is some $p \times 1$ vector of constants. The following example illustrates this condition:

Lemma 7 *Assume $H_0 : \psi = \psi_0$ where ψ_0 is the CF of the $\mathcal{N}(\mu, \sigma^2)$. Let f_0 be the pdf of the $\mathcal{N}(\mu, \sigma^2)$. The test T_B has only trivial power against local alternatives of the form*

$$H_{1n} : \psi_n(\tau) = \left(1 + \frac{ai\tau}{\sqrt{n}} - \frac{b\tau^2}{2\sqrt{n}}\right) \psi_0(\tau)$$

for some constants a and b . Moreover the density corresponding to ψ_n is

$$f_n(x) = \left\{1 + \frac{a}{\sqrt{n}} \frac{(x - \mu)}{\sigma^2} + \frac{b}{2\sqrt{n}} \left[\frac{(x - \mu)^2 - \sigma^2}{2\sigma^4}\right]\right\} f_0(x). \quad (19)$$

It follows from Lemma 7 that when μ and σ^2 are estimated, the test T_B has trivial power against alternatives of the form (19), which correspond to a second order Hermite expansion of the Gaussian density. The two additive terms in (19) contain the first two Hermite polynomials, which will be close to zero once μ and σ^2 are estimated. This is similar to what is found in other tests. For example, Bontemps and Meddahi (2005)'s moment test of normality cannot make use of the first two Hermite polynomials evaluated at the estimated parameters because their sample means will converge to 0 in probability even after scaling them by \sqrt{n} .

The following result establishes that T_B has power against all fixed alternatives (including those such that $L\eta = 0$).

Proposition 8 *Suppose Assumptions 1', 2-5 hold. The test T_B is consistent.*

In the next subsection, we analyze the power of our test in more detail.

4.2 Alternative with maximum power

We can follow the same steps as in Section 3.2 to characterize the alternatives for which the tests T_B have maximum power. Let $L^2(f_0)$ and assume the same normalization of u and the same power function $\Pi_B(a, c, u)$ as before. Following Neuhaus (1976) and Escanciano (2009), we can determine for which local alternative the test T_B has maximum power.

Let $h(x; t) = e^{itx} - \psi_0(\cdot; \theta_0)$. To analyze power, it is useful to rewrite the covariance operator \tilde{K} as $\tilde{T}^* \tilde{T}$, where \tilde{T} is an operator from $L^2(\pi)$ to $L^2(f_0)$ and \tilde{T}^* is the operator from $L^2(f_0)$ to $L^2(\pi)$, \tilde{T}^* being the adjoint of \tilde{T} . The operators \tilde{T} and \tilde{T}^* are as follows;

$$\begin{aligned} \tilde{T} &: L^2(\pi) \rightarrow L^2(f_0), \\ (\tilde{T}\varphi)(X) &= \int [\overline{h(X; t)} - \overline{D_0(t)} \langle BD_0, BD_0 \rangle^{-1} \overline{\langle B^* BD_0, h(X; \cdot) \rangle}] \varphi(t) \pi(t) dt, \\ \tilde{T}^* &: L^2(f_0) \rightarrow L^2(\pi), \text{ and} \\ (\tilde{T}^*\phi)(t) &= \int [h(x; t) - D_0(t) \langle BD_0, BD_0 \rangle^{-1} \langle B^* BD_0, h(x; \cdot) \rangle] \phi(x) f_0(x) dx. \end{aligned}$$

Moreover, $B\tilde{T}^*\tilde{T}B^* = B\tilde{K}B^*$ is compact and $\tilde{T}B^*$ admits a singular system $\{\tilde{a}_j, \tilde{\phi}_j, \tilde{\varphi}_j\}$ such that $\tilde{T}B^*\tilde{\phi}_j = \sqrt{\tilde{a}_j}\tilde{\varphi}_j$ and $B\tilde{T}^*\tilde{\varphi}_j = \sqrt{\tilde{a}_j}\tilde{\phi}_j$. $\tilde{\phi}_j$ are the eigenfunctions of $B\tilde{T}^*\tilde{T}B^*$ and $\tilde{\varphi}_j$ are the eigenfunctions of $\tilde{T}B^*B\tilde{T}^*$. The functions $\tilde{\varphi}_j$ can again be interpreted as principal components.

Observe that $L\eta = \tilde{T}^*u$. Hence,

$$\frac{\langle BL\eta, \tilde{\phi}_j \rangle}{\sqrt{\tilde{a}_j}} = \frac{\langle B\tilde{T}^*u, \tilde{\phi}_j \rangle}{\sqrt{\tilde{a}_j}} = \frac{\langle u, \tilde{T}B^*\tilde{\phi}_j \rangle}{\sqrt{\tilde{a}_j}} = \langle u, \tilde{\varphi}_j \rangle. \quad (20)$$

From (17) and Proposition 6, it follows that under $\tilde{H}_{un}(c)$,

$$T_B \xrightarrow{d} \sum_{j=1}^{\infty} \tilde{a}_j (e_j + c \langle u, \tilde{\varphi}_j \rangle)^2. \quad (21)$$

For those u such that $T^*u = 0$, the tests T_B have power equal to size. Therefore, we will focus on alternatives such that $T^*u \neq 0$, alternatives which belong to the orthogonal space to the null space of T^* (denoted $\mathcal{N}(T^*)$) – these are the alternatives corresponding to η such that $L\eta \neq 0$. For any compact operator T we have the relation, $\mathcal{N}(T^*)^\perp = \overline{\mathcal{R}(T)}$, where $\overline{\mathcal{R}(T)}$ is the closure of the range of T . Note that the sequence $\tilde{\varphi}_j$, for $j = 1, 2, \dots$ form a complete orthonormal basis of $\overline{\mathcal{R}(T)}$. Hence, the alternatives of interest are linear combinations of the $\tilde{\varphi}_j$. The analysis of the limiting distribution in (21) and the orthogonality of the $\tilde{\varphi}_j$ allow us to establish an analogous result to Proposition 4:

Proposition 9 *Suppose Assumptions 1', 2-5 hold. The limiting power $\Pi_B(\alpha, c, u)$ has the following properties.*

- (a) $\max_u \{\Pi_B(a, c, u) : u \in \overline{\mathcal{R}(T)}, \|u\|_{L^2(f_0)} = 1\} = \Pi_B(a, c, \tilde{\varphi}_1)$,
- (b) $\Pi_B(a, c, \tilde{\varphi}_j) \leq \Pi_B(a, c, \tilde{\varphi}_i)$ for $1 \leq i \leq j$,
- (c) $\lim_{j \rightarrow \infty} \Pi_B(a, c, \tilde{\varphi}_j) = a$.

As before, we observe that the maximum power is reached for the first principal component, and that power declines toward size a for subsequent directions.

5 J test when α goes to zero

5.1 Distribution under local alternatives

As we discussed at the end of Section 2, the continuum of moments analogue to the over-identification restrictions test diverges when α goes to zero, so we need to center and re-scale this statistic appropriately as in (15). But because \hat{q}_n in the denominator of this expression diverges as n goes to infinity, the rescaled test does not have power against contiguous alternatives. Therefore, we need to consider alternatives that converges to H_0 slower than the usual $n^{-1/2}$ rate. For that reason, in what follows we study the properties of J_{α_n} under local alternatives of the form

$$H_{2n} : \psi_n(t) = \psi_0(t) + b_n \eta(t)$$

where $\eta \in L^2(\pi)$, $\eta(0) = 0$, $\overline{\eta(t)} = \eta(-t)$, $|\eta(t)| < \mathfrak{C}$ for some constant \mathfrak{C} , and b_n is a sequence of numbers going to zero at a rate slower than \sqrt{n} . The precise rate will be specified later on. In this section, we assume for simplicity that the CF, ψ_0 , is completely specified under the null and the dependence in a known parameter, θ_0 , is omitted. The case where θ_0 is unknown is discussed in Remark 1 below. In the sequel, P_{2n} denotes the law of X_i under H_{2n} .

Under H_{2n} , we have

$$\sqrt{n}\{\hat{h}(\cdot) - E^{P_{2n}}[\hat{h}(\cdot)]\} \Rightarrow \mathcal{CN}(0, K, R)$$

in $L^2(\pi)$ where K and R are integral operators from $L^2(\pi)$ to $L^2(\pi)$ with kernels defined in (5) and (6).

Let $\{\lambda_j, \phi_j\}$ $i = 1, 2, \dots$ be the eigenvalues and eigenfunctions of K and $a_j = \frac{\lambda_j^2}{\lambda_j^2 + \alpha}$. Let $p_n = \sum_{j=1}^n a_j$, $q_n = \sum_{j=1}^n a_j^2$. Moreover, let \mathcal{H}_K be the reproducing kernel Hilbert space (RKHS) associated with K , defined as

$$\mathcal{H}_K = \left\{ f \in L^2(\pi) : \|f\|_K^2 = \sum \frac{\langle f, \phi_j \rangle^2}{\lambda_j} < \infty \right\}. \quad (22)$$

Assumption 6 $p_n/(q_n n \alpha) \rightarrow 0$ and $p_n^2/(q_n n) \rightarrow 0$ as n goes to infinity and α goes to zero.

Assumption 6 is very mild given that in Proposition 10 we will require $n\alpha^2 \rightarrow \infty$, and also from Lemma 9 in Carrasco and Florens (2000), it is known that if there exist $0 < \gamma < 1$ and some positive constant c such that $p_n \sim c\alpha^{-\gamma}$, then $q_n \sim e\alpha^{-\gamma}$ for some positive constant e (see also remark 2 below).

Proposition 10 *Suppose Assumptions 1', 4-6 hold. Assume that $\eta \in \mathcal{H}_K$ and*

$$\frac{nb_n^2}{\sqrt{q_n}} \rightarrow d \text{ for some constant } d. \quad (23)$$

Under H_{2n} , we have

$$J_{\alpha_n} \xrightarrow{d} \mathcal{N}(d \|\eta\|_K^2, 1)$$

as $n \rightarrow \infty$, $\alpha \rightarrow 0$, $n\alpha^2 \rightarrow \infty$, $nb_n\alpha^{3/4} \rightarrow \infty$, and $b_n^2\alpha^{-1} \rightarrow 0$, where $\|\cdot\|_K$ denotes the norm in the RKHS defined in (22).

Remark 1 Under H_0 , J_{α_n} converges to a standard normal distribution, therefore critical values from normal tables can be readily used. Moreover, when θ_0 is unknown and replaced by a consistent estimator, J_{α_n} converges again to a standard normal distribution under H_0 (the proof is similar to that of Proposition 10 and is omitted).

Remark 2 The condition (23) indicates the rate of b_n , which is related to the rate of the eigenvalues λ_j through q_n . Let us consider an example where $\lambda_j = j^{-m}$. Then $q_n \sim \alpha^{-1/(2m)}$ (see Carrasco and Florens (2000, Example 2) for the case $m = 1$ and Wahba (1975) for the general case). So condition (23) can be rewritten as $b_n \sim n^{-1/2}\alpha^{-1/(8m)}$. We observe that, in this case, the condition $n\alpha^2 \rightarrow \infty$ implies $nb_n\alpha^{3/4} \rightarrow \infty$ and $b_n^2/\alpha \rightarrow 0$.

Remark 3 The test J_{α_n} has nontrivial power against local alternatives $\psi_0 + cb_n\eta$ for any η .

Remark 4 The fact that J_{α_n} has trivial power against $1/\sqrt{n}$ alternatives is linked to the rescaling of the statistic. In fact, all tests involving centering and rescaling exhibit the same lack of power against contiguous alternatives. This includes Neyman's smooth test with an increasing number of polynomials (see Lehmann and Romano), the chi-square type test for conditional moments (De Jong and Bierens, 1994), the goodness-of-fit tests considered by Eubank and LaRiccia (1992), Härdle and Mammen (1993) and the one considered by Aït-Sahalia, Bickel, and Stocker (2001), among others.

Remark 5 Carrasco and Florens (2000) derived the asymptotic null distribution of J_{α_n} under a stronger assumption (Assumption 15: $q_n\sqrt{\alpha_n} \rightarrow \infty$). This assumption requires that the eigenvalues go to zero very slowly, which is not realistic here. On the contrary, the eigenvalues of K are likely to go to zero very fast, as illustrated in Figures 2e and 2f. For that reason, we propose a new proof which relaxes this assumption.

Remark 6 The lack of power of J_{α_n} against contiguous alternatives may speak in favor of tests such that T_B , which have power against contiguous alternatives. However, the test J_{α_n} may have higher power than T_B for higher frequency alternatives (case $j \rightarrow \infty$ in Proposition 8); see Theorem 3 in Eubank and LaRiccia (1992). The next remark considers this issue from a different angle.

Remark 7 Proposition 10 establishes the asymptotic distribution of J_{α_n} for η such that $\|\eta\|_K^2 < \infty$. However, this condition is not necessarily satisfied, so it is of special interest to look at what happens when it does not hold. Specifically, consider the case where

$$\frac{1}{\sqrt{q_n}} \sum_{l=1}^n \frac{a_l^2}{\lambda_l} |\langle \eta, \phi_l \rangle|^2 \rightarrow \infty. \quad (24)$$

The proof of Proposition 10 implies that the right rate for the alternatives H_{2n} is such that

$$nb_n^2 \frac{1}{\sqrt{q_n}} \sum_{l=1}^n \frac{a_l^2}{\lambda_l} |\langle \eta, \phi_l \rangle|^2 \rightarrow d$$

for some constant d . It follows from (24) that $nb_n^2 \rightarrow 0$. Hence the test J_{α_n} has power against local alternatives which approach the null hypothesis at a faster rate than $n^{-1/2}$. For these alternatives, the power of the tests T_B presented earlier remain $n^{-1/2}$. So the test J_{α_n} is able to detect certain alternatives which are closer to the null than the tests based on a fixed α . This result is similar to what was observed by Fan and Li (2000) in the context of specification tests for nonparametric regression. In particular, they show that nonparametric specification tests such as that of Härdle and Mammen (1993) with a fixed bandwidth has analogous properties as the integrated conditional tests of Bierens (1982) and Bierens and Ploberger (1997). Further, they show that kernel based tests with bandwidth going to zero can detect specific alternatives at a faster rate than $n^{-1/2}$. As we mentioned before, we can interpret α as a bandwidth in our tests.

5.2 Numerical invariance to moment transformations

As is well known, the traditional J test corresponding to the continuously updated estimator (CUE) is invariant to parameter-dependent linear transformations of the moments (see Hansen, Heaton and Yaron (1995)). To illustrate this fact, let $\hat{\underline{h}}(\theta)$ be the sample average of a vector of moments and M_θ be a (possibly complex-valued) square invertible matrix. Then, it is easy to check that the J-test based on $\hat{\underline{h}}(\theta)$ is the same as the J-test based on $M_\theta \hat{\underline{h}}(\theta)$ because

$$J = n \hat{\underline{h}}(\theta)^* M_\theta^* (M_\theta \hat{\mathcal{K}}_\theta M_\theta^*)^{-1} M_\theta \hat{\underline{h}}(\theta) = n \hat{\underline{h}}(\theta)^* \hat{\mathcal{K}}_\theta^{-1} \hat{\underline{h}}(\theta).$$

When one uses regularization to invert the covariance matrix, this result is not true in general. Indeed, we have that

$$n \hat{\underline{h}}(\theta)^* M_\theta^* (M_\theta \hat{\mathcal{K}}_\theta M_\theta^*)^{1/2} [(M_\theta \hat{\mathcal{K}}_\theta M_\theta^*)^2 + \alpha I]^{-1} (M_\theta \hat{\mathcal{K}}_\theta M_\theta^*)^{1/2} M_\theta \hat{\underline{h}}(\theta)$$

is not usually equal to

$$n\hat{h}(\theta)^* \hat{\mathcal{K}}_\theta^{1/2} (\hat{\mathcal{K}}_\theta^2 + \alpha I)^{-1} \hat{\mathcal{K}}_\theta^{1/2} \hat{h}(\theta)$$

unless M_θ is unitary, that is $M_\theta M_\theta^* = M_\theta^* M_\theta = I$, in which case the two expressions coincide.

When there is a continuum of moment conditions, an analogous result turns out to be true for unitary transformations of h .

Define U_θ as a nonrandom linear operator from $L^2(\pi)$ into $L^2(\pi)$. Let U_θ^* be the adjoint of U_θ . By the Riesz representation theorem, there is a unique $g_\theta(\cdot, s)$ such that

$$(U_\theta \varphi)(s) = \langle g_\theta(\cdot, s), \varphi(\cdot) \rangle = \int g_\theta(t, s) \varphi(t) \pi(t) dt.$$

Let K_θ be the covariance operator of $h_i(\cdot; \theta)$ and \tilde{K}_θ be the covariance operator of $U_\theta h_i(\cdot; \theta)$. The kernel of \tilde{K}_θ is such that

$$\begin{aligned} \tilde{k}_\theta(s_1, s_2) &= E[(U_\theta h_i(\cdot; \theta))(s_1) (U_\theta h_i(\cdot; \theta))(s_2)^*] \\ &= E \left[\int g_\theta(t, s_1) h_i(t; \theta) \pi(t) dt \int \overline{g_\theta(u, s_2) h_i(u; \theta) \pi(u) du} \right] \\ &= \int g_\theta(t, s_1) \left\{ \int E[h_i(t; \theta) \overline{h_i(u; \theta)}] \overline{g_\theta(u, s_2) \pi(u)} du \right\} \pi(t) dt \\ &= \langle g_\theta(\cdot, s_1), K_\theta g_\theta(\cdot, s_2) \rangle. \end{aligned}$$

Then, we can characterize \tilde{K}_θ :

$$\begin{aligned} (\tilde{K}_\theta \varphi)(\tau) &= \int \int g_\theta(t, \tau) \left\{ \int E[h_i(t; \theta) \overline{h_i(u; \theta)}] \overline{g_\theta(u, s) \pi(u)} du \right\} \pi(t) dt \varphi(s) \pi(s) ds \\ &= (U_\theta K_\theta U_\theta^* \varphi)(\tau). \end{aligned}$$

Proposition 11 *Let U_θ be an unitary operator from $L^2(\pi)$ to $L^2(\pi)$ i.e. $U_\theta^* U_\theta = U_\theta U_\theta^* = I$.*

Then, the following equality holds:

$$\left\| U_\theta \hat{h}(\theta) \right\|_{(U_\theta K_\theta U_\theta^*)^\alpha} = \left\| \hat{h}(\theta) \right\|_{K_\theta^\alpha} \quad (25)$$

regardless of the sample size n .

This means that the CUE versions of tests T_B with $B = (K^\alpha)^{-1/2}$ (α fixed) and $J_{\alpha_n}(\hat{\theta}, \hat{K})$ are numerically invariant to unitary transformations of h . For non unitary transformations, the result is no longer true because of the regularization. In contrast, T_B with $B = I$ for instance is not even invariant to unitary transformations.

6 Monte Carlo experiments

In this section, we assess the finite sample performance of our proposed tests by means of several extensive Monte Carlo exercises. In addition, we compare them to several popular nonparametric tests based on the empirical distribution function, as well as to directional tests that target specific parametric alternatives to the null. In all cases, our sample size is $n = 100$.

6.1 Testing univariate normality

The first design we consider is a univariate normal distribution, which is by far the most common null hypothesis in distributional tests. In order to make our tests numerically invariant to affine transformations of the observations, we systematically centre and standardize them using the sample mean and standard deviation (with denominator n), which are the ML estimators under the null. As proved by Carrasco and Florens (2014), an asymptotically equivalent procedure would estimate the mean and variance by minimizing the continuum of moment conditions criterion function, but this would result in an increase of the computational costs. Either way, we can set the true mean and variance to 0 and 1, respectively, without loss of generality.

We consider three versions of our test, which differ in the way the covariance operator is estimated. The first one uses the theoretical covariance operator for a standard normal, which we presented in Section 2. In turn, the second and third versions rely on the centred and uncentred sample estimators using expressions (12) and (13), respectively, with the matrices C and \hat{C} computed using the analytical integrals in Appendix A.2. Given that these two sample versions produce very similar results, we only report the centred one in what follows. Importantly, the test that uses the theoretical covariance operator offers two notable computational advantages: i) the calculation of its eigenvalues and eigenfunctions depends on the number of grid points M , which we set to 1,000, but not on the sample size, so it can be used with very large datasets; and ii) we only need to compute those eigenelements once regardless of the number of Monte Carlo simulations.

In view of the discussion in Section 2, we look at two values of the Tikhonov regularization parameter α (.1 and .01) and two values for the scale parameter of the $\mathcal{N}(0, \omega^2)$ density defining inner products (1 and $\sqrt{10}$). As we have previously discussed, increasing ω not only changes the eigenvalues and eigenfunctions, but more intuitively, it pays relative more attention to the characteristic function for large (in absolute terms) values of its argument t .

In this univariate context, it is straightforward to compute the Cramer von Mises (CvM), Kolmogorov-Smirnov (KS) and Anderson-Darling (AD) statistics on the basis of the probability integral transforms (PIT) of the standardized observations obtained through the standard nor-

mal cdf (see Appendix A.3 for details). Their usual asymptotic distributions are invalid, though, because those PITs make use of the sample mean and variance.

Further, we also compute two moment-based tests: one focusing on the fourth Hermite polynomial $(z^4 - 3z^2 + 1)/\sqrt{24}$ and another one that simultaneously looks at the third Hermite polynomial $(z^3 - 3z)/\sqrt{6}$ too. The advantage of working with Hermite polynomials is that they are asymptotically invariant to parameter estimation under the null (see e.g. Bontemps and Meddahi (2005)). As is well known, these two statistics can be derived as Lagrange multiplier tests against a variety of non-normal distributions (see e.g. Jarque and Bera (1980) or Mencía and Sentana (2012)). Finally, we also compute the Bierens and Wang (2012) test described in Appendix A.3 using a MATLAB translation of their C+ code.

The first thing we do is to compute all the aforementioned tests for 10,000 simulated samples generated under the null, whence we obtain finite sample critical values. This parametric bootstrap procedure automatically generates size-adjusted rejection rates, as forcefully argued by Horowitz and Savin (2000); see also Dufour (2006) for a discussion of Monte Carlo tests.

Panels A-F of Table 1 contain those rejection rates for six different alternatives: a symmetric Student t with 12 degrees of freedom; an asymmetric Student t with the same number of degrees of freedom but skewness parameter $\beta = -.75$; a scale mixture of two normals with the same kurtosis as the symmetric t , 3.75, and mixture probability $\lambda = .1$ (outlier case); another scale mixture with the same kurtosis but $\lambda = .75$ (inlier case); a location-scale mixture constructed in such a way that it has same skewness and kurtosis as the normal and $E(x^5) = -1$, $E(x^6) = 18$; and finally, the second order Hermite expansion of the normal density mentioned in Lemma 7 with parameters $a = .4$ and $b = .5$. Details on how we simulate those distributions can be found in Supplemental Appendix A.4. Figure B3 in the Supplemental Appendix presents the densities of all the alternative distributions once they have been standardized so that they all have 0 means and unit standard deviations in the population.

The first four columns of each panel in Table 1 report the results for the test that is based on the theoretical covariance operator, $J(\hat{\theta}, K_{\hat{\theta}})$, for the different values of α and ω that we consider. In turn, the next four columns contain the same figures for the test $J(\hat{\theta}, \hat{K})$ which uses the centred sample estimator of the covariance operator instead. As can be seen across the different panels, in all cases the results seem robust to the choice of the regularization parameter α . For the majority of the DGPs, $J(\hat{\theta}, K_{\hat{\theta}})$ has more power when $\omega = 1$ while the performance of $J(\hat{\theta}, \hat{K})$ is better with $\omega = \sqrt{10}$. In addition, they generally outperform the other consistent tests that we consider, with AD being the most powerful of them. Somewhat surprisingly, this is also true when the DGP is the second order Hermite expansion of the normal mentioned

in Lemma 7 (Panel F). Nevertheless, it is important to remember that this lemma refers to local alternatives, while our test is consistent versus fixed alternatives. Not surprisingly, the LM tests are the most powerful testing procedures when the distribution under the alternative is the one they are designed to detect. Specifically, this applies to S- t , the LM test against symmetric Student t alternatives, in Panel A, and A- t , the LM test against asymmetric Student t alternatives, in Panel B.

In summary, our proposed tests display good power against a variety of alternatives.

6.2 Testing uniformity

The second design we consider is a uniform distribution. Although this distribution does not often arise as a model for natural phenomena, it plays a fundamental role in statistics for two reasons: most computer-based pseudo-random number generators aim to draw uniform variates, and the PITs of any continuous random variables are uniform. To facilitate the comparison with the normal distribution, we transform the standard uniform random numbers by subtracting from them their population mean (.5) and scaling them up by their population standard deviation ($\sqrt{12}$), so that the resulting distribution will become standardized.

We consider exactly the same versions of our tests as in Section 6.1, but with the expressions for the population kernel and the centred and uncentered sample versions modified accordingly, as explained in Appendix A.2. We also compute the three non-parametric tests based on the CDF, as well as the Bierens and Wang (2012) test. As for directional tests, we consider two possibilities. The first one is the LM test of uniform vs beta proposed by Sefton (1992), which exploits the fact that a beta distribution with shape parameters $a = b = 1$ becomes uniform. This test is based on the average scores with respect to the beta parameters evaluated under the null, which are $1 + \ln(u)$ and $1 + \ln(1 - u)$, respectively.² The second directional test is a moment test based on the first two Jacobi polynomials evaluated again under the null, namely $\sqrt{3}(2u - 1)$ and $\sqrt{3}(6u^2 - 6u + 1)$, which was proposed by Bontemps and Meddahi (2012). As is well known, those polynomials constitute an orthonormal basis for the beta random variable.

The three panels of Table 2 contain the parametric bootstrap rejection rates for three different alternatives. The first one is a symmetric, unimodal beta distribution with parameters $a = b = 1.1$. The second one is an asymmetric unimodal concave beta distribution with parameters $a = 1.1$ and $b = 1$. Finally, the last distribution is generated as the standard Gaussian PITs of observations drawn from the same asymmetric Student t distribution with 12 degrees of freedom

²The asymptotic variance for the scores reported by Sefton (1992) seems to be incorrect. For that reason, we use instead 1 for the two asymptotic variances and $(6 - \pi^2)/6$ for the covariance. Therefore, the LM test is $T/2$ times the square of the difference between the two scores divided by $\pi^2/6$ plus the square of their sum divided by $(12 - \pi^2/6)$.

and asymmetric parameter as in Section 6.1. The motivation for including this alternative is that we can use it to compare the direct application of our proposed tests to the original observations and to a monotonic transformation of them. Figure B4 in the Supplemental Appendix presents the densities of these alternative distributions.

The first four columns of each panel report the results for the test that based on the theoretical covariance operator, $J(\theta_0, K)$, for the same values of α and the scale parameter ω of the $\mathcal{N}(0, \omega^2)$ density defining inner products as in the previous subsection, while the next four columns focus on $J(\theta_0, \hat{K})$. As in Section 6.1, the rejection rates of our tests seem robust to the choice of the regularization parameter α . But in this case they are also less sensitive to the choice of ω . As before, the test based on K outperforms the one that uses centred sample estimator of the covariance operator \hat{K} . Interestingly, both of them outperform the competitors when the DGP is either a symmetric beta or the Gaussian PITs of observations drawn from an asymmetric Student t . In contrast, CvM and AD are slightly more powerful when the alternative is the asymmetric beta. Somewhat surprisingly, the LM test is not particularly powerful.

6.3 Testing bivariate normality

Our next design is a bivariate normal distribution, which is by far the most common null hypothesis in multivariate distributional tests. Once again, we make our tests numerically invariant to affine transformations of the observations by systematically centring and standardizing them using the sample mean and the Cholesky decomposition of the sample covariance matrix (with denominator n), which are the ML estimators under the null.³ Thus, we can set the true means and standard deviations to 0 and 1, respectively, and the correlation coefficient to 0 without loss of generality.

We consider exactly the same versions of our tests as in the Section 6.1, but with the expressions for the population kernel and the centred and uncentred sample versions modified accordingly (see Appendix A.2). However, we do not compute any classical non-parametric tests because there is no consensus on distribution-free multivariate generalization of the CvM, KS and AD statistics based on the joint distribution function. Nevertheless, we continue to apply the Bierens and Wang (2012) test. By analogy with the univariate normal case in section 6.1, we also consider two directional tests: the LM test of a multivariate normal against a multivariate Student t in Fiorentini, Sentana and Calzolari (2003) (denoted S- t), which effectively focuses on Mardia's (1970) coefficient of multivariate excess kurtosis, and the LM test against a generalized

³As we mentioned before, an asymptotically equivalent procedure would estimate the two means and variances as well as the covariance by minimizing the continuum of moment conditions criterion function, but this would result in a huge increase of the computational cost.

hyperbolic distribution in Mencía and Sentana (2012) (denoted A- t), which also looks at third moments in order to capture asymmetries in the multivariate distribution. By construction, both tests are asymptotically invariant to parameter estimation under the null.

The three panels of Table 3 contain the parametric bootstrap rejection rates for three different alternatives. The first one is a multivariate Student t with 12 degrees of freedom. The second one is an asymmetric Student t with the same degrees of freedom and vector of asymmetric parameters $(-.75, -.75)$. Finally, the third alternative is a spherically symmetric bivariate version of the outlier distribution considered in Section 6.1. Figure B5 in the Supplemental Appendix presents the densities of these alternative distributions.

As in Table 1, the first four columns of each panel in Table 3 report the results for the test $J(\hat{\theta}, K_{\hat{\theta}})$ again for the same values of α and the scale parameter ω of the $\mathcal{N}(0, \omega^2)$ density defining inner products as in subsection 6.1 and the next four columns correspond to same figures for the test $J(\hat{\theta}, \hat{K})$. As can be seen in Table 3, in all cases the results seem robust to the choice of the regularization parameter α . Moreover, for the DGPs we consider $J(\hat{\theta}, K_{\hat{\theta}})$ has more power when $\omega = 1$ while the performance of $J(\hat{\theta}, \hat{K})$ is better with $\omega = \sqrt{10}$, as in the univariate case. Interestingly, $J(\hat{\theta}, K_{\hat{\theta}})$ beats the S- t LM test when the DGPs is asymmetric Student t and there is a tie between $J(\hat{\theta}, \hat{K})$ and S- t LM test when the alternative is a discrete-scale mixture of normals.

6.4 Testing chi-square

Another design we consider is a chi-square distribution with two degrees of freedom. Like the uniform, the chi-square distribution does not often arise as a model for natural phenomena. But it also plays a fundamental role in statistics because it is the distribution of the (square) Mahalanobis distance of a multivariate normal random variable from its mean. In other words, it corresponds to the distribution of $(y_i - \mu)' \Sigma^{-1} (y_i - \mu)$ when $y_i \sim \mathcal{N}(\mu, \Sigma)$.

We consider exactly the same versions of our tests as in Sections 6.1 and 6.2, but with the expressions for the population kernel and the centred and uncentred sample versions in Appendix A.2 suitably modified. In that regard, the main difference is that we define inner products using a uniform density over $[-\omega, \omega]$, for values of ω equal to 1 and $\sqrt{10}$. Although we standardize again the random draws by subtracting their population mean (=2) and scaling them down by their population standard deviation (=2), their distribution remains asymmetric, which implies that both the CF and the eigenfunctions of the associated covariance operator are complex, as explained in Section 2. This creates a normalization problem because any complex vector of unit length remains so after scaling its elements by any complex scalar on the unit

circle, e^{iv} , where $v \in [0, 2\pi)$. Nevertheless, our proposed tests are numerically invariant to any chosen normalization.

We also compute the three non-parametric tests, as well as the Bierens and Wang (2012) test. As for directional tests, we consider two possibilities. The first one is the LM test of chi square with N degrees of freedom versus F with the same number of degrees of freedom in the numerator but ν degrees of freedom in the denominator proposed by Fiorentini, Sentana and Calzolari (2003). This test is based on the average score with respect to the reciprocal of ν evaluated under the null, which coincides with the second order Laguerre polynomial

$$\frac{1}{4}\zeta^2 - 2\zeta + 2,$$

whose asymptotic variance for $N = 2$ is 4 under the null. The second directional test is the LM test against a gamma distribution with mean N but shape parameter $\alpha \neq N/2$ developed in Amengual and Sentana (2012). In this case, the score is proportional to

$$\left(\frac{\zeta}{2} - 1\right) - \left[\ln\left(\frac{\zeta}{2}\right) - \psi(1)\right],$$

whose asymptotic variance is $\psi'(1) - 1$, where $\psi(\cdot)$ and $\psi'(\cdot)$ are the digamma and trigamma functions, respectively.

The three panels of Table 4 contain the parametric bootstrap rejection rates for three different alternatives. The first one is an F distribution with 12 degrees of freedom in the denominator, while the second one is a gamma distribution with shape parameter $\alpha = 2/3$ and scale parameter $\beta = 3$. Finally, the last distribution is generated as the square norm of observations drawn from a bivariate asymmetric Student t distribution with 12 degrees of freedom. Once again, the motivation for including this alternative is that we can use it to compare the direct application of our proposed bivariate Gaussian tests to the original observations or to a transformation of them which implicitly imposes spherical symmetry. In that regard, the F distribution would correspond to a bivariate Student t while the gamma to a Kotz distribution. The densities of these alternative distributions are reported in Figure B6 in the Supplemental Appendix.

As in Table 2, the first four columns of each panel of Table 4 report the results for the test $J(\theta_0, K)$, for the different values of α and ω that we consider, while the next four columns contain the same figures for $J(\theta_0, \hat{K})$. Once again, the results seem robust to the choice of the regularization parameter α , but at the same time they are less sensitive to the choice of ω . Still, for $J(\theta_0, K)$ the value $\omega = 1$ delivers higher rejection rates. As before, the test based on the theoretical covariance operator outperforms the one using centred sample estimator of the covariance operator. Interestingly, $J(\theta_0, K)$ has more power than its competitors, except when

the DGP is Gamma.

6.5 Testing Cauchy

The last design we consider is a Cauchy distribution with location and scale parameters 0 and 1, respectively. In order to make our tests numerically invariant to affine transformations of the observations, we systematically centre and standardize them using the ML estimators of location and scale under the null.

We consider exactly the same versions of our tests as in Section 6.1, but with the expressions for the population kernel and the centred and uncentered sample versions modified accordingly, as explained in Appendix A.2. We also compute the three non-parametric tests based on the CDF, as well as the Bierens and Wang (2012) test.

The three panels of Table 4 contain the parametric bootstrap rejection rates for three different alternatives. The first one is a Student t with 2 degrees of freedom, while in Panel B we draw from an asymmetric Student t with 6 degrees of freedom and skewness parameter $\beta = -.25$. Finally, the last distribution we consider is a Laplace with location and scale parameters 0 and $1/\sqrt{2}$, respectively. Details on how we simulate those distributions can be found in Appendix A.4. Figure B7 in the Supplemental Appendix presents the densities of these alternative distributions.

The first four columns of each panel in Table 1 report the results for the test that is based on the theoretical covariance operator, $J(\hat{\theta}, K_{\hat{\theta}})$, once again for the different values of α and ω that we consider. In turn, the next four columns contain the same figures for the test $J(\hat{\theta}, \hat{K})$ which uses the centred sample estimator of the covariance operator. As can be seen across the different panels, in all cases the results seem robust to the choice of the regularization parameter α . For the majority of the DGPs, both $J(\hat{\theta}, K_{\hat{\theta}})$ and $J(\hat{\theta}, \hat{K})$ have more power when $\omega = \sqrt{10}$. In addition, they generally outperform the other consistent tests that we consider, with BW being the most powerful among them.

Once again, our proposed tests display good power against a variety of alternatives.

7 Conclusion

In this paper we propose goodness-of-fit tests based on comparing the empirical and theoretical characteristic functions. Our proposals are based on the continuum of moment conditions analogue to the usual overidentifying restrictions test, and therefore take into account the correlation between the influence functions for different argument values.

We consider different versions depending on whether the parameter vector θ is known in advance or replaced by a consistent estimator, and whether we make use of the analytical

expression for the covariance operator or estimate it. Relying on the theoretical covariance operator offers substantial computational gains because the calculation of its eigenvalues and eigenvectors does not depend on the sample size, which allows its use with very large datasets.

We derive the asymptotic distribution of our proposed tests for fixed regularization parameter and when this vanishes with the sample size. Both types of tests have very different asymptotic properties. The fixed α J test has a nonstandard asymptotic distribution which depends on nuisance parameters but has power against $1/\sqrt{n}$ alternatives. In contrast, the vanishing α J test has a standard normal asymptotic distribution but generally fails to reject local $1/\sqrt{n}$ alternatives, except for some specific alternatives which it can detect at a faster rate.

Our theoretical study of power sheds some light on the alternatives for which each test is more powerful. While there is no test whose power dominates overall, it seems that fixing α at a small positive value is a good compromise. An extensive Monte Carlo exercise confirms this point by showing that our proposed tests display good power in finite samples against a variety of alternatives.

Although we have focused on a random sample framework for pedagogical reasons, versions of our tests robust to serial or cross-sectional dependence in the observations should be relatively straightforward. The analysis of conditional distributions would also constitute a very valuable but non-trivial addition with many potentially interesting empirical applications.

Appendix

A Proofs and auxiliary results

Proof of Lemma 1. K is self-adjoint positive definite because it is a covariance operator ($k(s, t) = \overline{k(t, s)}$) and its null space is reduced to 0, i.e. $Kf = 0 \Rightarrow f = 0$ (see the proof of Proposition A.1, condition A.5(i) in Carrasco et al (2007)). K is a Hilbert-Schmidt operator because its kernel is square integrable, indeed

$$\int \int |k(s, t)|^2 \pi(s) ds \pi(t) dt < \infty.$$

Consequently, K admits an infinite spectrum of positive eigenvalues. Let $\{\lambda_j, \varphi_j\}$ be the eigenvalues arranged in decreasing order and eigenfunctions (the eigenfunctions are taken orthonormal in $L^2(\pi)$) of K . By Mercer's formula (see Carrasco, Florens, and Renault, 2007, Theorem 2.42),

$$k(t, s) = \sum_j \lambda_j \varphi_j(t) \varphi_j(s).$$

By setting $s = t$, we have

$$\sum \lambda_j = \int k(t, t) \pi(t) dt.$$

Here $k(t, s) = \psi(t-s) - \psi(t)\psi(-s)$. Hence $k(t, t) = 1 - |\psi(t)|^2 \leq 1$. It follows that $\sum \lambda_j \leq 1$, which in turn implies that $0 \leq \lambda_j \leq 1$ because the operator is self-adjoint positive definite. Therefore $\lambda_j^2 \leq \lambda_j$ and hence $\sum \lambda_j^2 \leq 1$. So the Hilbert-Schmidt norm of K is also bounded by 1:

$$\|K\|_{HS}^2 = \int \int |k(t, s)|^2 \pi(s) ds \pi(t) dt = \sum \lambda_j^2 \leq 1,$$

as desired. □

Proof of Proposition 2. We check the conditions (a) to (c) of Lemma 3.1 of Chen and White (1998)⁴ on

$$W_{nj} = \frac{1}{\sqrt{n}} \left(h_j - \frac{c\eta}{\sqrt{n}} \right).$$

Checking (a): We need to check that for all $\varphi \in L^2(\pi)$, $\sum_{j=1}^n \langle W_{nj}, \varphi \rangle \xrightarrow{d} \mathcal{CN}(0, \sigma^2(\varphi), \delta(\varphi))$ where $\sigma^2(\varphi) = \langle \varphi, K\varphi \rangle > 0$ and $\delta(\varphi) = \langle \varphi, R\varphi \rangle$. To do so, first notice that under H_{1n} , $W_{nj} = \frac{1}{\sqrt{n}} [e^{itX_j} - \psi_n(t)]$. We have $E[\langle W_{nj}, \varphi \rangle] = 0$ and $\langle W_{nj}, \varphi \rangle$, $j = 1, 2, \dots, n$ are indepen-

⁴The results of Chen and White (1998) are stated for real random variables, but we adapt them here to complex variables.

dent. Moreover,

$$\begin{aligned}
E[|\langle W_{nj}, \varphi \rangle|^2] &= E[\langle W_{nj}, \varphi \rangle \overline{\langle W_{nj}, \varphi \rangle}] \\
&= E \int \int W_{nj}(s) \overline{\varphi(s) W_{nj}(t) \varphi(t)} \pi(s) ds \pi(t) dt \\
&= \int \int E[W_{nj}(s) \overline{W_{nj}(t)}] \overline{\varphi(s)} \varphi(t) \pi(s) ds \pi(t) dt \\
&= \frac{1}{n} \langle \varphi, K_n \varphi \rangle
\end{aligned}$$

where K_n is the integral operator with kernel

$$\begin{aligned}
&k_n(s, t) \\
&= \psi_n(s-t) - \psi_n(s) \psi_n(-t) \\
&= \psi_0(s-t) - \psi_0(s) \psi_0(-t) + \frac{c\eta(s-t)}{\sqrt{n}} - \frac{c\eta(s)}{\sqrt{n}} \psi_0(-t) - c\psi_0(s) \frac{\eta(-t)}{\sqrt{n}} + c^2 \frac{\eta(s)\eta(-t)}{n}.
\end{aligned}$$

Interchanging the order of integration is justified by the fact that $\frac{1}{n} \langle \varphi, K_n \varphi \rangle < \infty$. Now, we check the conditions of Lindeberg-Feller central limit theorem (van der Vaart (1998), Proposition 2.27) to establish $\sum_{j=1}^n \langle W_{nj}, \varphi \rangle \xrightarrow{d} \mathcal{CN}(0, \sigma^2(\varphi), \delta(\varphi))$. Let $Y_{nj} = \langle W_{nj}, \varphi \rangle$. Here Y_{nj} are independent scalar random variables with zero mean and finite variance. The three conditions for the CLT are

$$\begin{aligned}
\text{(i)} \quad &\sum_{j=1}^n E[|Y_{nj}|^2 I\{|Y_{nj}| > \varepsilon\}] \rightarrow 0 \text{ for every } \varepsilon > 0, \\
\text{(ii)} \quad &\sum_{j=1}^n E(Y_{nj} \overline{Y_{nj}}) \rightarrow \sigma^2(\varphi), \text{ and} \\
\text{(iii)} \quad &\sum_{j=1}^n E(Y_{nj}^2) \rightarrow \delta(\varphi).
\end{aligned}$$

Note that

$$\begin{aligned}
|Y_{nj}|^2 &= |\langle W_{nj}, \varphi \rangle|^2 \\
&\leq \left\| \frac{1}{\sqrt{n}} [e^{itx_j} - \psi_n(t)] \right\|^2 \|\varphi\|^2 \\
&\leq \frac{C}{n} \|\varphi\|^2
\end{aligned}$$

for some fixed constant C . Hence,

$$\begin{aligned} \sum_{j=1}^n E[|Y_{nj}|^2 I\{|Y_{nj}| > \varepsilon\}] &\leq \frac{C \|\varphi\|^2}{n} \sum_{j=1}^n P\{|Y_{nj}| > \varepsilon\} \\ &\leq \frac{C \|\varphi\|^2}{n} \sum_{j=1}^n \frac{E[|Y_{nj}|^2]}{\varepsilon^2} \\ &\leq \frac{C^2}{n\varepsilon^2} \|\varphi\|^4 \end{aligned}$$

by Markov inequality. So condition (i) is satisfied. For (ii), we use the results above which give

$$\sum_{j=1}^n E(Y_{nj} \overline{Y_{nj}}) = \sum_{j=1}^n E[|\langle W_{nj}, \varphi \rangle|^2] = \langle \varphi, K_n \varphi \rangle \rightarrow \langle \varphi, K \varphi \rangle,$$

and hence, (ii) is also satisfied. Finally,

$$\sum_{j=1}^n E(Y_{nj}^2) = \sum_{j=1}^n E[\langle W_{nj}, \varphi \rangle^2] = \langle \varphi, R_n \varphi \rangle \rightarrow \langle \varphi, R \varphi \rangle$$

where R_n is the integral operator with kernel $r_n(s, t) = k_n(s, -t)$, hence (iii) follows.

Checking (b) and (c): By Remark 3.3 (ii) of Chen and White (1998), conditions (b) and (c) can be replaced by the following condition:

W_{nj} is strictly stationary and

$$\lim_{n \rightarrow \infty} E \left| \sum_{j=1}^n W_{nj} \right|^2 \leq C < \infty. \quad (\text{A1})$$

We have

$$\begin{aligned} E \left| \sum_{j=1}^n W_{nj} \right|^2 &= E \left\langle \sum_{j=1}^n W_{nj}, \sum_{l=1}^n \overline{W_{nl}} \right\rangle \\ &= \sum_{j=1}^n E \langle W_{nj}, \overline{W_{nj}} \rangle \\ &= \sum_{j=1}^n k_n(s, s) \\ &= 1 - |\psi_0(s)|^2 - \frac{c\eta(s)}{\sqrt{n}} \psi_0(-s) - \psi_0(s) \frac{c\eta(-s)}{\sqrt{n}} + \frac{c^2 |\eta(s)|^2}{n} \end{aligned}$$

which is bounded because $|\psi_0(s)|^2 \leq 1$ by the property of CFs and $\eta(s)$ is also bounded.

Therefore, (A1) is satisfied and $\sum_{j=1}^n W_{nj}$ is tight.

It follows that $\sqrt{n}\hat{h} = \sum_{j=1}^n W_{nj} + c\eta \Rightarrow \mathcal{CN}(c\eta, K, R)$. \square

The two lemmas below will be used in the proof of Proposition 3.

Lemma 12 *Let U be a complex Gaussian process in $L^2(\pi)$ such that*

$$U \sim \mathcal{CN}(c\eta, K, R)$$

where K and R are arbitrary covariance and relation operators.

(a) *Let B be a bounded operator, then $Z \equiv BU \sim \mathcal{CN}(cB\eta, BKB^*, BR\bar{B}^*)$ where $\bar{B}^* = \overline{B^*}$ is the adjoint of the complex conjugate of B (or equivalently the complex conjugate of the adjoint of B).*

(b) *Let $\underline{Z} = (Z, \bar{Z})'$. Then, \underline{Z} is a complex bivariate Gaussian process with mean $c(B\eta, \bar{B}\eta)'$ and covariance operator Γ defined as the operator from $(L^2(\pi))^2$ to $(L^2(\pi))^2$ such that*

$$g = \begin{pmatrix} g_1 \\ g_2 \end{pmatrix} \in (L^2(\pi))^2 \rightarrow \Gamma g = \begin{pmatrix} BKB^* & BR\bar{B}^* \\ \bar{B}\bar{R}B^* & \bar{B}\bar{K}\bar{B}^* \end{pmatrix} g.$$

(c) *Let $Y = (\text{Re}(Z), \text{Im}(Z))$. Then Y is a real bivariate Gaussian process with mean $c(\text{Re}(B\eta), \text{Im}(B\eta))' \equiv c\tilde{\eta}$ and covariance operator $M\Gamma M^*$. Moreover,*

$$\|BU\|^2 = \|Y\|_2^2 \sim \sum_{j=1}^{\infty} b_j \chi_j^2(1, \delta_j) + c \sum_l \langle \tilde{\eta}, \zeta_l \rangle_2^2$$

where b_j, ζ_j are the nonzero eigenvalues and eigenfunctions of $M\Gamma M^*$, ζ_l are the eigenfunctions of $M\Gamma M^*$ associated with the zero eigenvalue and $\delta_j = c^2 \langle \tilde{\eta}, \zeta_j \rangle_2^2 / b_j$.

Proof of Lemma 12. (a) and (b) can be proved by direct algebra. As for (c), we first show that $\|BU\|^2 = \|Z\|^2 = \|Y\|_2^2$. Indeed, we have

$$\begin{aligned} \|Z\|^2 &= \langle Z, Z \rangle \\ &= \int (\text{Re}(Z(t)) + i\text{Im}(Z(t))) (\text{Re}(Z(t)) - i\text{Im}(Z(t))) \pi(t) dt \\ &= \langle \text{Re}(Z), \text{Re}(Z) \rangle + \langle \text{Im}(Z), \text{Im}(Z) \rangle \\ &= \|Y\|_2^2. \end{aligned}$$

Note that Y is a real process which satisfies $Y = M\underline{Z}$, so its mean is $cM(B\eta, \bar{B}\eta)' = c(\text{Re}(B\eta), \text{Im}(B\eta))' \equiv c\tilde{\eta}$ and its variance is $M\Gamma M^*$. Let V denote a bivariate Gaussian process $\mathcal{N}(0, M\Gamma M^*)$. We have

$$\|Y\|_2^2 \stackrel{d}{\sim} \|c\tilde{\eta} + V\|_2^2 = \sum_{i=j,l} (c \langle \tilde{\eta}, \zeta_i \rangle_2 + \langle V, \zeta_i \rangle_2)^2$$

where the equality uses the Karhunen-Loeve theorem of Gaussian processes and the fact that the 2×1 eigenfunctions ζ_i ($i = j, l$) form an orthonormal basis of $(L^2(\pi))^2$. Moreover,

$$V = \sum_{i=1}^{\infty} \langle V, \zeta_j \rangle_2 \zeta_j = \sum_{j=1}^{\infty} \sqrt{b_j} \frac{\langle V, \zeta_j \rangle_2}{\sqrt{b_j}} \zeta_j$$

where $\langle V, \zeta_j \rangle_2 / \sqrt{b_j}$ are *iid* $\mathcal{N}(0, 1)$ and $\langle V, \zeta_l \rangle_2 = 0$ because its mean=0 and its variance=0. So the result follows. \square

Lemma 13 *Consider the case where K and R have kernels defined in (5) and (6) respectively. When π is symmetric around 0 and $B = I$ or $B = (K^\alpha)^{-1/2}$, then the distribution of $\|Y\|^2$ in Lemma 12 simplifies to $\sum_{j=1}^{\infty} a_j \chi_j^2(1, \delta_j)$ where $\delta_j = c^2 \langle B\eta, \phi_j \rangle^2 / a_j$, a_j and ϕ_j are the eigenvalues and eigenfunctions of BKB^* .*

Proof of Lemma 13. First consider the case $B = I$. Let us compute the spectrum of Γ . Let ϕ_j be the orthonormal eigenfunctions of K associated with the eigenvalues λ_j . Let $\tilde{\phi}_j(s) = \phi_j(-s)$. It turns out that for every ϕ_j of K , there are two 2×1 eigenfunctions of Γ , namely $\tilde{\zeta}_j(s) = (\phi_j(s), \tilde{\phi}_j(s))'$ associated with the eigenvalues $2\lambda_j$ and $\tilde{\zeta}_l(s) = (\phi_l(s), -\tilde{\phi}_l(s))'$ (with $l = j$) associated with the eigenvalue 0. Moreover, the eigenfunctions are orthogonal, so that we get a complete eigenvalue-eigenfunctions decomposition for the Gaussian process considered. Indeed

$$\begin{aligned} \left(\Gamma \tilde{\zeta}_j \right) (s) &= \begin{pmatrix} \int k(s, t) \phi_j(t) \pi(t) dt + \int k(s, -t) \phi_j(-t) \pi(t) dt \\ \int k(-s, t) \phi_j(t) \pi(t) dt + \int k(-s, -t) \phi_j(-t) \pi(t) dt \end{pmatrix} \\ &= \begin{pmatrix} \int k(s, t) \phi_j(t) \pi(t) dt + \int k(s, u) \phi_j(u) \pi(u) du \\ \int k(-s, t) \phi_j(t) \pi(t) dt + \int k(-s, u) \phi_j(u) \pi(u) du \end{pmatrix} \\ &= \begin{pmatrix} 2\lambda_j \phi_j(s) \\ 2\lambda_j \phi_j(-s) \end{pmatrix} \\ &= 2\lambda_j \tilde{\zeta}_j(s) \end{aligned}$$

by the change of variable $u = -t$ and using the fact that $\pi(-u) = \pi(u)$. Moreover, $\left(\Gamma \tilde{\zeta}_l \right) (s) = 0$. We see that Γ is singular. The eigenfunctions of $M\Gamma M^*$ are $\zeta_j = M\tilde{\zeta}_j$ (associated with the eigenvalues $a_j = \lambda_j$) and $\zeta_l = M\tilde{\zeta}_l$ (associated with the 0 eigenvalue). Indeed, using the fact that $M^*M = \frac{1}{2}I$,

$$\begin{aligned} M\Gamma M^* M\tilde{\zeta}_j &= \frac{1}{2} M\Gamma \tilde{\zeta}_j \\ &= \lambda_j M\tilde{\zeta}_j. \end{aligned}$$

Because MTM^* is a real matrix, its eigenfunctions ζ_j have to be real. The ϕ_j are not defined uniquely because we can multiply complex eigenfunctions by a complex number on the unit circle without altering their norm. We select ϕ_j so that ζ_j is real. The fact that it is possible to transform ϕ_j into $e^{id}\phi_j$ (for some constant d) so that ζ_j is real, is proved in Lemma 14 below. Now, we compute $\langle \tilde{\eta}, \zeta_i \rangle$:

$$\begin{aligned} \langle \tilde{\eta}, \zeta_i \rangle &= \langle \tilde{\eta}, M\tilde{\zeta}_j \rangle \\ &= \langle M^*\tilde{\eta}, \tilde{\zeta}_j \rangle \\ &= \frac{1}{2} \left\{ \langle \eta, \phi_j \rangle + \langle \bar{\eta}, \tilde{\phi}_j \rangle \right\} \\ &= \langle \eta, \phi_j \rangle \end{aligned}$$

using the fact that $\bar{\eta}(t) = \eta(-t)$ and π is symmetric around 0. Moreover,

$$\begin{aligned} \langle \tilde{\eta}, \zeta_l \rangle &= \langle M^*\tilde{\eta}, \tilde{\zeta}_l \rangle \\ &= \frac{1}{2} \left\{ \langle \eta, \phi_l \rangle - \langle \bar{\eta}, \tilde{\phi}_l \rangle \right\} \\ &= 0. \end{aligned}$$

When $B = (K^\alpha)^{-1/2}$, the proof is similar as above because B has the same eigenfunctions as K . Details are omitted. \square

Lemma 14 *Let φ_j be an orthonormal eigenfunction of K , then there exists a constant d so that $\phi_j(t) = e^{id}\varphi_j(t)$, $\tilde{\phi}_j(t) = \phi_j(-t)$, $\tilde{\zeta}_j(t) = \left(\phi_j(t), \tilde{\phi}_j(t) \right)'$ and $\tilde{\zeta}_l(t) = \left(\phi_l(t), -\tilde{\phi}_l(t) \right)'$ are such that $\zeta_j = M\tilde{\zeta}_j$ and $\zeta_l = M\tilde{\zeta}_l$ are real. Moreover, ϕ_j is such that $\langle \hat{h}, \phi_j \rangle$ is real for all $j = 1, 2, \dots$*

Proof of Lemma 14. Let us denote $\text{Re}(\varphi_j(t)) = a_t$ and $\text{Im}(\varphi_j(t)) = b_t$. Here, we treat the case $a_t \neq 0$ and $b_t \neq 0$ (the cases where either a_t or b_t is null can be treated similarly). We have

$$\begin{aligned} \phi_j &= e^{id}\varphi_j = (\cos d + i \sin d)(a_t + ib_t) \\ &= a_t \cos d - b_t \sin d + i(a_t \sin d + b_t \cos d) \end{aligned}$$

and

$$\zeta_j = M\tilde{\zeta}_j = \frac{1}{2} \begin{pmatrix} \phi_j + \tilde{\phi}_j \\ i(\tilde{\phi}_j - \phi_j) \end{pmatrix}.$$

Hence, ζ_j is real if and only if $\text{Im}(\phi_j + \tilde{\phi}_j) = 0$ and $\text{Re}(\tilde{\phi}_j - \phi_j) = 0$, which is equivalent to

$$(a_t + a_{-t}) \sin d + (b_t + b_{-t}) \cos d = 0, \quad (\text{A2})$$

$$(a_{-t} - a_t) \cos d - (b_{-t} - b_t) \sin d = 0 \quad (\text{A3})$$

or equivalently

$$\tan d = \frac{-(b_t + b_{-t})}{a_t + a_{-t}} = \frac{(a_t - a_{-t})}{(b_t - b_{-t})}.$$

For this to be possible, we need $a_t^2 + b_t^2 = a_{-t}^2 + b_{-t}^2$ which is equivalent to

$$|\varphi_j(t)|^2 = |\varphi_j(-t)|^2. \quad (\text{A4})$$

Now we show that (A4) holds for any eigenfunction φ_j of the covariance operator K . As K is a compact self-adjoint operator, we have

$$\begin{aligned} k(s, s) &= \sum_j \lambda_j |\varphi_j(s)|^2, \\ k(s, s) &= k(-s, -s) = \sum_j \lambda_j |\varphi_j(-s)|^2. \end{aligned}$$

It follows that $\sum_j \lambda_j [|\varphi_j(s)|^2 - |\varphi_j(-s)|^2] = 0$ for all s . Hence $|\varphi_j(s)|^2 = |\varphi_j(-s)|^2$. The case of $\tilde{\zeta}_l$ associated with 0 can be treated in the same manner and we get the same condition (A4).

Using (A2) and (A3), one can check that ϕ_j satisfies the relation $\phi_j(t) = \overline{\phi_j(-t)}$. Consequently, $\langle \hat{h}, \phi_j \rangle$ is real for all $j = 1, 2, \dots$ because

$$\begin{aligned} \langle \hat{h}, \phi_j \rangle &= \int \hat{h}(t) \overline{\phi_j(t)} \pi(t) dt \\ &= \int \hat{h}(t) \phi_j(-t) \pi(t) dt \\ &= \int \hat{h}(-s) \phi_j(s) \pi(s) ds \\ &= \int \overline{\hat{h}(s)} \phi_j(s) \pi(s) ds \\ &= \overline{\langle \hat{h}, \phi_j \rangle} \end{aligned}$$

using a change of variable, $s = -t$, and the property that $\overline{\hat{h}(s)} = \hat{h}(-s)$. \square

Proof of Proposition 3. Adapting the results of Chen and White (1992, working paper) to complex processes and taking into account that B is bounded, we have

$$B_n \sqrt{n} \hat{h} \Rightarrow \mathcal{CN} \left(cB\eta, BKB^*, BR\overline{B}^* \right)$$

where B^* is the adjoint of B . Then the results follow from Lemma 12 and Lemma 13. \square

Proof of Proposition 4. The proof is similar to those of Neuhaus (1976, Theorem 2.2.) and Escanciano (2009, Theorem 1) and is not repeated here for brevity. \square

Proof of Proposition 5. Under our assumptions,

$$\left\| B_n \hat{h}(\cdot; \theta) \right\| \xrightarrow{P_0} \left\| BE^{P_0}[h_j(\cdot; \theta)] \right\|$$

uniformly in θ . (The uniformity part comes from the fact that $\hat{h}(\cdot; \theta) - E[h_j(\cdot; \theta)] = \frac{1}{n} \sum_{j=1}^n e^{itX_j} - \psi_0(t; \theta_0)$ does not depend on θ .) Moreover, $E[h_j(\cdot; \theta)] = \psi_0(\cdot; \theta_0) - \psi_0(\cdot; \theta)$. By the identification assumption, the objective function reaches its minimum at $\theta = \theta_0$. Hence, $\hat{\theta}$ is consistent under H_0 .

We turn our attention toward the asymptotic normality. To simplify the notation, we write $\psi_0(\theta)$ for $\psi_0(\cdot; \theta)$ and $\hat{h}(\theta)$ for $\hat{h}(\cdot; \theta)$, and $\frac{\partial \psi_0(\hat{\theta})}{\partial \theta}$ for $\frac{\partial \psi_0(\theta)}{\partial \theta} \Big|_{\theta=\hat{\theta}}$. The first order condition of the minimization problem gives

$$\begin{aligned} & \left\langle B_n \frac{\partial \psi_0(\hat{\theta})}{\partial \theta}, B_n \hat{h}(\hat{\theta}) \right\rangle \\ = 0 & = \left\langle B_n \frac{\partial \psi_0(\hat{\theta})}{\partial \theta}, B_n \hat{h}(\theta_0) \right\rangle - \left\langle B_n \frac{\partial \psi_0(\hat{\theta})}{\partial \theta}, B_n \frac{\partial \psi_0(\tilde{\theta})}{\partial \theta} (\hat{\theta} - \theta_0) \right\rangle \end{aligned}$$

where $\tilde{\theta}$ is between θ_0 and $\hat{\theta}$. It follows that

$$\sqrt{n}(\hat{\theta} - \theta_0) = \left\langle B_n \frac{\partial \psi_0(\hat{\theta})}{\partial \theta}, B_n \frac{\partial \psi_0(\tilde{\theta})}{\partial \theta} \right\rangle^{-1} \left\langle B_n \frac{\partial \psi_0(\hat{\theta})}{\partial \theta}, B_n \sqrt{n} \hat{h}(\theta_0) \right\rangle$$

By the continuity of $\frac{\partial \psi_0}{\partial \theta}$ and the consistency of $\hat{\theta}$, we have

$$\sqrt{n}(\hat{\theta} - \theta_0) = \langle BD_0, BD_0 \rangle^{-1} \left\langle B^* BD_0, \sqrt{n} \hat{h}(\cdot; \theta_0) \right\rangle + o_{P_0}(1). \quad (\text{A5})$$

The asymptotic normality follows from Proposition 1.

For the convergence of $\hat{\theta}$ to θ_1 under H_1 , we use the same arguments as for the consistency under H_0 . The existence of the minimum comes from the fact that $\psi_0(\cdot; \theta)$ is continuous in θ and Θ is compact. \square

Proof of Proposition 6.

(i) The mean value theorem gives

$$\begin{aligned}
\sqrt{n}\hat{h}(\hat{\theta}) &= \sqrt{n}\hat{h}(\theta_0) - \frac{\partial\psi_0(\hat{\theta})}{\partial\theta}\sqrt{n}(\hat{\theta} - \theta_0) \\
&= \sqrt{n}\hat{h}(\theta_0) - D_0\sqrt{n}(\hat{\theta} - \theta_0) + o_{P_0}(1) \\
&= \sqrt{n}\hat{h}(\theta_0) - D_0\langle BD_0, BD_0\rangle^{-1}\left\langle B^*BD_0, \sqrt{n}\hat{h}(\theta_0)\right\rangle + o_{P_0}(1)
\end{aligned}$$

by equation (A5). In turn, the contiguity of P_n to P_0 implies that

$$\sqrt{n}\hat{h}(\hat{\theta}) - \sqrt{n}\hat{h}(\theta_0) + D_0\langle BD_0, BD_0\rangle^{-1}\left\langle B^*BD_0, \sqrt{n}\hat{h}(\theta_0)\right\rangle \xrightarrow{P_n} 0. \quad (\text{A6})$$

By Proposition 2, we have under H_{1n}

$$\sqrt{n}\hat{h}(\theta_0) - D_0\langle BD_0, BD_0\rangle^{-1}\left\langle B^*BD_0, \sqrt{n}\hat{h}(\theta_0)\right\rangle \Rightarrow \mathcal{CN}(L\eta, \tilde{K}, \tilde{R}) \quad (\text{A7})$$

where \tilde{R} is the relation operator whose explicit expression is not given because it is not needed. Combining Equations (A6) and (A7) yields $\sqrt{n}\hat{h}(\hat{\theta}) \Rightarrow \mathcal{CN}(L\eta, \tilde{K}, \tilde{R})$ under H_{1n} . The kernel of \tilde{K} can be computed explicitly as follows:

$$\begin{aligned}
&\tilde{k}(s, t) \\
&= E\left[(\sqrt{n}\hat{h}(s) - D_0(s)\langle BD_0, BD_0\rangle^{-1}\left\langle B^*BD_0, \sqrt{n}\hat{h}\right\rangle)\right. \\
&\quad \left.\times(\sqrt{n}\hat{h}(t) - D_0(t)\langle BD_0, BD_0\rangle^{-1}\left\langle B^*BD_0, \sqrt{n}\hat{h}\right\rangle)\right].
\end{aligned}$$

The detailed calculation for one of the four terms gives

$$\begin{aligned}
&E\left[D_0(s)\langle BD_0, BD_0\rangle^{-1}\left\langle B^*BD_0, \sqrt{n}\hat{h}\right\rangle\sqrt{n}\hat{h}(t)\right] \\
&= D_0(s)\langle BD_0, BD_0\rangle^{-1}E\left[\int B^*BD_0(u)\sqrt{n}\hat{h}(u)\pi(u)du\sqrt{n}\hat{h}(t)\right] \\
&= D_0(s)\langle BD_0, BD_0\rangle^{-1}\int B^*BD_0(u)E[h_j(u)\overline{h_j(t)}]\pi(u)du \\
&= D_0(s)\langle BD_0, BD_0\rangle^{-1}(KB^*BD_0)(t).
\end{aligned}$$

The other terms can be computed similarly.

(ii) The proof of (ii) is similar to that of Proposition 3 and hence omitted. \square

Proof of Lemma 7. The CF of a $\mathcal{N}(\mu, \sigma^2)$ is $\psi_0(t) = e^{i\mu t - \frac{\sigma^2 t^2}{2}}$. Let $\theta = (\mu, \sigma^2)'$, then

$$D_0 = \frac{\partial\psi_0}{\partial\theta} = \begin{pmatrix} it\psi_0(t) \\ -\frac{t^2}{2}\psi_0(t) \end{pmatrix}.$$

Let $v = (a, b)$ and $\eta(t) = v'D_0 = \left(ait - \frac{bt^2}{2}\right)\psi_0(t)$. Now consider $\psi_n(t) = \left(1 + \frac{ait}{\sqrt{n}} - \frac{bt^2}{2\sqrt{n}}\right)\psi_0(t)$.

Observe that $\psi_n(0) = 1$, $\overline{\psi_n(t)} = \psi_n(-t)$. We need $|\psi_n(t)| < 1$ which will be satisfied if $b > 0$ (and possibly for $b < 0$ and n large enough). So ψ_n satisfies the necessary conditions to be a CF, however these conditions are not sufficient. Necessary and sufficient conditions for a function ψ_n to be CF are that (a) $\psi_n(0) = 1$, and (b) ψ_n is non-negative definite (see Theorem 4.2.2 of Lukacs (1960)). It can be shown that, given ψ_0 is a CF, ψ_n will satisfy (b) for n large enough. So ψ_n is a CF.

Moreover, $\psi_n(t)$ is absolutely integrable so the density (f_n) corresponding to ψ_n satisfies:

$$\begin{aligned} f_n(x) &= \frac{1}{2\pi} \int e^{-itx} \psi_n(t) dt \\ &= \frac{1}{2\pi} \int e^{-itx} \left(1 + \frac{ait}{\sqrt{n}} - \frac{bt^2}{2\sqrt{n}} \right) \psi_0(t) dt \\ &= \frac{1}{2\pi} \int e^{-itx} \psi_0(t) dt + \frac{ai}{\sqrt{n}} \frac{1}{2\pi} \int t e^{-itx} \psi_0(t) dt - \frac{b}{2\sqrt{n}} \frac{1}{2\pi} \int e^{-itx} t^2 \psi_0(t) dt. \end{aligned}$$

Note that

$$\begin{aligned} \frac{i}{2\pi} \int t e^{-itx} \psi_0(t) dt &= \frac{\partial \frac{1}{2\pi} \int e^{-itx} \psi_0(t) dt}{\partial \mu}, \\ \frac{1}{2\pi} \int e^{-itx} t^2 \psi_0(t) dt &= -2 \frac{\partial \frac{1}{2\pi} \int e^{-itx} \psi_0(t) dt}{\partial \sigma^2}. \end{aligned}$$

At the same time,

$$\frac{1}{2\pi} \int e^{-itx} \psi_0(t) dt = \frac{1}{\sqrt{2\pi\sigma^2}} \exp \left[-\frac{(x-\mu)^2}{2\sigma^2} \right] \equiv f_0(x),$$

$$\begin{aligned} \frac{\partial f_0(x)}{\partial \sigma^2} &= -\frac{1}{2} \frac{1}{\sqrt{2\pi\sigma^3}} \exp \left[-\frac{(x-\mu)^2}{2\sigma^2} \right] \\ &\quad + \frac{(x-\mu)^2}{2\sigma^4} \frac{1}{\sqrt{2\pi\sigma^2}} \exp \left[-\frac{(x-\mu)^2}{2\sigma^2} \right] \\ &= \left[\frac{(x-\mu)^2}{2\sigma^4} - \frac{1}{2\sigma^2} \right] f_0(x), \end{aligned}$$

$$\frac{\partial f_0(x)}{\partial \mu} = \frac{(x-\mu)}{\sigma^2} f_0(x).$$

It then follows that $f_n(x) = \left\{ 1 + \frac{a}{\sqrt{n}} \frac{(x-\mu)}{\sigma^2} + \frac{b}{2\sqrt{n}} \left[\frac{(x-\mu)^2 - \sigma^2}{2\sigma^4} \right] \right\} f_0(x)$. \square

Proof of Proposition 8. Under H_1 , $\hat{h}(\cdot; \hat{\theta}) \xrightarrow{P_1} E^{P_1} h_j(\cdot; \theta_1) = \psi(\cdot) - \psi_0(\cdot; \theta_1) \neq 0$, where $\psi(\cdot)$ is the CF of X_j under H_1 , and then the result follows. \square

Preliminary results to the proof of Proposition 10.

The following lemmas will be used in the proof of Proposition 10.

Let $Y_{in}(s)$ be the process defined as

$$Y_{in}(s) = e^{isX_i} - \psi_n(s).$$

Under H_{2n} , Y_{in} $i = 1, 2, \dots$ are *iid* with mean 0 and covariance $E[Y_{in}(s) \overline{Y_{in}(t)}] = \psi_n(s-t) - \psi_n(s)\psi_n(-t) \equiv k_n(s, t)$. Let K_n be the integral operator with kernel k_n and $(\lambda_{l,n}, \phi_{l,n})$ be the eigenvalues and eigenfunctions of K_n . Note that K_n converges to K when n goes to infinity.

Lemma 15 Under H_{2n} , $\left(\frac{\langle Y_{in}, \phi_{l,n} \rangle}{\sqrt{\lambda_{l,n}}}\right)$, $l = 1, 2, \dots$ are uncorrelated across l with zero mean and variance equal to 1.

Proof of Lemma 15. We have

$$\begin{aligned} E \left[\langle Y_{in}, \phi_{l,n} \rangle \overline{\langle Y_{in}, \phi_{l',n} \rangle} \right] &= E \int Y_{in}(s) \overline{\phi_{l,n}(s)} \pi(s) ds \int \overline{Y_{in}(t)} \phi_{l',n}(t) \pi(t) dt \\ &= \int \overline{\phi_{l,n}(s)} \int E \left[Y_{in}(s) \overline{Y_{in}(t)} \right] \phi_{l',n}(t) \pi(t) dt \pi(s) ds \\ &= \overline{\langle \phi_{l,n}, K_n \phi_{l',n} \rangle} \\ &= \begin{cases} \lambda_{l,n} & \text{if } l = l', \\ 0 & \text{otherwise,} \end{cases} \end{aligned}$$

as desired. □

The following lemma is taken from Eubank and LaRiccia (1992) and is reproduced here for convenience. Note that in our setting, Y_{in} is complex but we can still apply this lemma because w_{ijn} is real.

Lemma 16 (Lemma 2 of Eubank and LaRiccia (1992)) Let $\{Y_{in}\}_{i=1}^n$, $n = 1, 2, \dots$ be a triangular array of random variables that are *iid* within rows. Set $w_{ijn} = w_{ijn}(Y_{in}, Y_{jn}) + w_{ijn}(Y_{jn}, Y_{in})$ for some function $w_{ijn}(\cdot, \cdot)$ and assume that $E[w_{ijn}|Y_{in}] = 0$ for all $i, j \leq n$. Define

$$\begin{aligned} w(n) &= \sum_{1 \leq i < j \leq n} w_{ijn}, \\ \sigma(n)^2 &= \text{Var}(w(n)) = \sum_{1 \leq i < j \leq n} E(w_{ijn}^2), \\ G_I &= \sum_{1 \leq i < j \leq n} E(w_{ijn}^4), \\ G_{II} &= \sum_{1 \leq i < j < k \leq n} [E(w_{ijn}^2 w_{ikn}^2) + E(w_{jin}^2 w_{jkn}^2) + E(w_{kin}^2 w_{kjn}^2)], \end{aligned}$$

and

$$G_{IV} = \sum_{1 \leq i < j < k < m \leq n} [E(w_{ijn}w_{ikn}w_{mjn}w_{mkn}) + E(w_{ijn}w_{imn}w_{kjn}w_{kmn}) + E(w_{imn}w_{ikn}w_{jkn}w_{jmn})].$$

Then, if G_I , G_{II} , and G_{IV} are all of smaller order than $\sigma(n)^4$,

$$\frac{w(n)}{\sigma(n)} \xrightarrow{d} \mathcal{N}(0, 1).$$

Lemma 17 Let $a_{l,n} = \frac{\lambda_{l,n}^2}{\lambda_{l,n}^2 + \alpha}$, $p_{n,n} = \sum_{j=1}^n a_{l,n}$, $q_{n,n} = 2 \sum_{j=1}^n a_{l,n}^2$. Under H_{2n} :

$$\frac{\sum_{l=1}^n \frac{a_{l,n}}{\lambda_{l,n}} \left| \left\langle \sqrt{n} \hat{h}, \phi_{l,n} \right\rangle \right|^2 - p_{n,n}}{\sqrt{q_{n,n}}} \xrightarrow{d} \mathcal{N}\left(c \|\eta\|_K^2, 1\right)$$

as $n \rightarrow \infty$, $\alpha \rightarrow 0$, $p_{n,n}^2/(q_{n,n}n) \rightarrow 0$, and $p_{n,n}/(q_{n,n}n\alpha) \rightarrow 0$.

Proof of Lemma 17. Our proof draws from the proof of Theorem 1 in Eubank and LaRiccia (1992). Here and in the subsequent proofs of results for Proposition 10, all the expectations are computed under H_{2n} . Dropping the subscript n from $a_{l,n}$, $\lambda_{l,n}$, $\phi_{l,n}$, $p_{n,n}$, and $q_{n,n}$, we obtain

$$\begin{aligned} & \frac{\sum_{l=1}^n \frac{a_l}{\lambda_l} \left| \left\langle \sqrt{n} \hat{h}, \phi_l \right\rangle \right|^2 - p_n}{\sqrt{q_n}} \\ &= \frac{\sum_{l=1}^n \frac{a_l}{\lambda_l} \left| \left\langle \sqrt{n} \{ \hat{h} - E(\hat{h}) + E(\hat{h}) \}, \phi_l \right\rangle \right|^2 - p_n}{\sqrt{q_n}} \\ &= \frac{\sum_{l=1}^n \frac{a_l}{\lambda_l} \left| \left\langle \sqrt{n} \{ \hat{h} - E(\hat{h}) \}, \phi_l \right\rangle \right|^2 - p_n}{\sqrt{q_n}} + R_n \end{aligned}$$

where

$$\begin{aligned} R_n &= \frac{2 \operatorname{Re} \sum_{l=1}^n \frac{a_l}{\lambda_l} \left\langle \sqrt{n} \{ \hat{h} - E(\hat{h}) \}, \phi_l \right\rangle \overline{\left\langle \sqrt{n} E(\hat{h}), \phi_l \right\rangle}}{\sqrt{q_n}} \\ &+ \frac{\sum_{l=1}^n \frac{a_l}{\lambda_l} \left| \left\langle \sqrt{n} E(\hat{h}), \phi_l \right\rangle \right|^2}{\sqrt{q_n}}. \end{aligned}$$

In a first step, we will show that R_n converges to $d \|\eta\|_K^2$ in probability under H_{2n} as n goes to infinity and α goes to zero. In a second step, we will show that, under H_{2n} ,

$$\frac{\sum_{l=1}^n \frac{a_l}{\lambda_l} \left| \left\langle \sqrt{n} \{ \hat{h} - E(\hat{h}) \}, \phi_l \right\rangle \right|^2 - p_n}{\sqrt{q_n}} \xrightarrow{d} \mathcal{N}(0, 1).$$

First step. We have $E(\hat{h}) = \psi_n - \psi_0 = b_n \eta$ and

$$E(R_n) = \frac{nb_n^2}{\sqrt{q_n}} \sum_{l=1}^n \frac{a_l}{\lambda_l} \langle \eta, \phi_l \rangle^2.$$

Moreover, $\sum_{l=1}^n \frac{a_l}{\lambda_l} |\langle \eta, \phi_l \rangle|^2 \rightarrow \sum_{l=1}^{\infty} \frac{1}{\lambda_l} |\langle \eta, \phi_l \rangle|^2 = \|\eta\|_K^2$ and $\frac{nb_n^2}{\sqrt{q_n}} \rightarrow d$ as n goes to infinity and α goes to zero. Therefore, $E(R_n) \rightarrow d \|\eta\|_K^2$.

Now we show that the variance of R_n goes to zero. Using the notation $Y_i = Y_{in}$, we have

$$V(R_n) = V \left[\frac{2 \operatorname{Re} \sum_{l=1}^n \frac{a_l}{\lambda_l} \left\langle \frac{\sum_{i=1}^n Y_i}{\sqrt{n}}, \phi_l \right\rangle \overline{\langle \sqrt{nb_n} \eta, \phi_l \rangle}}{\sqrt{q_n}} \right].$$

Using the fact that $V(\operatorname{Re}(Z)) \leq V(Z)$ for any complex random variable Z , we have

$$\begin{aligned} V(R_n) &\leq \frac{4nb_n^2}{q_n} V \left[\frac{\sum_{i=1}^n \sum_{l=1}^n \frac{a_l}{\lambda_l} \langle Y_i, \phi_l \rangle \overline{\langle \eta, \phi_l \rangle}}{\sqrt{n}} \right] \\ &= \frac{4nb_n^2}{q_n} V \left[\sum_{l=1}^n \frac{a_l}{\lambda_l} \langle Y_i, \phi_l \rangle \overline{\langle \eta, \phi_l \rangle} \right] \end{aligned}$$

because $Y_i, i = 1, 2, \dots, n$ are *iid*. As $\langle Y_i, \phi_l \rangle, l = 1, 2, \dots$ are uncorrelated by Lemma 15, we obtain

$$\begin{aligned} V(R_n) &\leq \frac{4nb_n^2}{q_n} \sum_{l=1}^n \frac{a_l^2}{\lambda_l} |\langle \eta, \phi_l \rangle|^2 \\ &\leq \frac{4nb_n^2}{q_n} \|\eta\|_K^2 \rightarrow 0. \end{aligned}$$

It follows that R_n converges to $d \|\eta\|_K^2$ in probability under H_{2n} .

Second step. We have

$$\begin{aligned} \frac{\sum_{l=1}^n \frac{a_l}{\lambda_l} \left| \left\langle \sqrt{n} \{ \hat{h} - E(\hat{h}) \}, \phi_l \right\rangle \right|^2 - p_n}{\sqrt{q_n}} &= \frac{\sum_{l=1}^n \frac{a_l}{\lambda_l} \left| \left\langle \frac{1}{\sqrt{n}} \sum_{i=1}^n Y_i, \phi_l \right\rangle \right|^2 - p_n}{\sqrt{q_n}} \\ &= \frac{w_1(n) + w(n)}{\sqrt{q_n}} \end{aligned}$$

where

$$\begin{aligned} w_1(n) &= \frac{1}{n} \sum_{l=1}^n \sum_{i=1}^n \frac{a_l}{\lambda_l} |\langle Y_i, \phi_l \rangle|^2 - p_n, \\ w(n) &= \frac{2}{n} \sum_{l=1}^n \frac{a_l}{\lambda_l} \sum_{1 \leq i < j \leq n} \operatorname{Re} \langle Y_i, \phi_l \rangle \overline{\langle Y_j, \phi_l \rangle} = \sum_{1 \leq i < j \leq n} w_{ijn}, \end{aligned}$$

with

$$w_{ijn} = \frac{2}{n} \sum_{l=1}^n \frac{a_l}{\lambda_l} \operatorname{Re} \langle Y_i, \phi_l \rangle \overline{\langle Y_j, \phi_l \rangle}.$$

First, we show that $w_1(n)/\sqrt{q_n} \xrightarrow{P_{2n}} 0$. We have

$$\begin{aligned} E[w_1(n)] &= \frac{1}{n} \sum_{l=1}^n \sum_{i=1}^n \frac{a_l}{\lambda_l} E[|\langle Y_i, \phi_l \rangle|^2] - p_n \\ &= \sum_{l=1}^n a_l - p_n = 0. \end{aligned}$$

As $\langle Y_i, \phi_l \rangle^2$ are independent across i , we have

$$\begin{aligned} &V[w_1(n)] \\ &= \frac{1}{n} V \left[\sum_{l=1}^n a_l \left(\frac{|\langle Y_i, \phi_l \rangle|^2}{\lambda_l} - 1 \right) \right] \\ &= \sum_{l=1}^n \frac{a_l^2}{n} E \left[\left(\frac{|\langle Y_i, \phi_l \rangle|^2}{\lambda_l} - 1 \right)^2 \right] + \sum_{l \neq l'} \frac{a_l a_{l'}}{n} E \left[\left(\frac{|\langle Y_i, \phi_l \rangle|^2}{\lambda_l} - 1 \right) \left(\frac{|\langle Y_i, \phi_{l'} \rangle|^2}{\lambda_{l'}} - 1 \right) \right] \end{aligned} \quad (\text{A8})$$

Using Lemma 15, we have

$$\begin{aligned} &\sum_{l \neq l'} \frac{a_l a_{l'}}{n} E \left[\left(\frac{|\langle Y_i, \phi_l \rangle|^2}{\lambda_l} - 1 \right) \left(\frac{|\langle Y_i, \phi_{l'} \rangle|^2}{\lambda_{l'}} - 1 \right) \right] \\ &= \sum_{l \neq l'} \frac{a_l a_{l'}}{n} E \left(\frac{|\langle Y_i, \phi_l \rangle|^2}{\lambda_l} \frac{|\langle Y_i, \phi_{l'} \rangle|^2}{\lambda_{l'}} \right) \end{aligned} \quad (\text{A9})$$

$$- \sum_{l \neq l'} \frac{a_l a_{l'}}{n}. \quad (\text{A10})$$

Consider (A10): We have

$$\frac{1}{q_n} \sum_{l \neq l'} \frac{a_l a_{l'}}{n} \leq \frac{p_n^2}{q_n n},$$

which goes to zero by assumption. To deal with the term (A9), we exploit the fact that for n large enough, $|Y_i| = |e^{itX_i} - \psi_n(t)| \leq |e^{itX_i}| + |\psi_n(t)| = 2$, hence $\|Y_i\|^2 \leq 4$ and $|\langle Y_i, \phi_l \rangle|^2 \leq 4$ by Cauchy-Schwarz and $\|\phi_l\| = 1$. Therefore, by Lemma 15,

$$E \left(\frac{|\langle Y_i, \phi_l \rangle|^2}{\lambda_l} \frac{|\langle Y_i, \phi_{l'} \rangle|^2}{\lambda_{l'}} \right) \leq \frac{4}{\lambda_l \lambda_{l'}} E \left(|\langle Y_i, \phi_{l'} \rangle|^2 \right) = \frac{4}{\lambda_l}.$$

Hence,

$$\sum_{l \neq l'} \frac{a_l a_{l'}}{n} E \left(\frac{|\langle Y_i, \phi_l \rangle|^2}{\lambda_l} \frac{|\langle Y_i, \phi_{l'} \rangle|^2}{\lambda_{l'}} \right) \leq \sum_{l \neq l'} \frac{a_l a_{l'}}{n \lambda_l} = \frac{p_n}{n} \sum_l \frac{a_l}{\lambda_l}.$$

Note that

$$\sum_l \frac{a_l}{\lambda_l} = \sum_l \frac{\lambda_l}{\lambda_l^2 + \alpha} \leq \frac{1}{\alpha} \sum_l \lambda_l.$$

So, we obtain:

$$\frac{p_n}{q_n n} \sum_l \frac{a_l}{\lambda_l} \leq \frac{p_n}{q_n n \alpha}$$

which goes to zero by assumption.

The first term in (A8) can be treated in the same manner. Thus, $V[w_1(n)]/q_n \rightarrow 0$ under our assumptions and hence $w_1(n)/\sqrt{q_n} \xrightarrow{P_{2n}} 0$.

Second, we show that under H_{2n}

$$\frac{w(n)}{\sqrt{q_n}} \xrightarrow{d} \mathcal{N}(0, 1).$$

To establish this result, we check all the conditions of Lemma 16.

$$\sigma(n)^2 = V(w_n) = \sum_{1 \leq i < j \leq n} E(w_{ij}^2),$$

where, using the fact that $|\operatorname{Re} Z| \leq |Z|$ for all complex Z , we have

$$\begin{aligned} E(w_{ij}^2) &\leq \frac{4}{n^2} E \left[\left(\sum_{l=1}^n \frac{a_l}{\lambda_l} |\langle Y_i, \phi_l \rangle| |\langle Y_j, \phi_l \rangle| \right)^2 \right] \\ &= \frac{4}{n^2} \sum_{l=1}^n \frac{a_l^2}{\lambda_l^2} E[|\langle Y_i, \phi_l \rangle|^2 |\langle Y_j, \phi_l \rangle|^2] \\ &= \frac{4}{n^2} \sum_{l=1}^n \frac{a_l^2}{\lambda_l^2} E[|\langle Y_i, \phi_l \rangle|^2] E[|\langle Y_j, \phi_l \rangle|^2] \\ &= \frac{4}{n^2} \sum_{l=1}^n a_l^2 = \frac{2q_n}{n^2} \end{aligned}$$

because the $\langle Y_i, \phi_l \rangle$ are uncorrelated across l and independent across i . Hence,

$$\sigma(n)^2 \sim q_n.$$

Consider now the term G_I :

$$G_I = \sum_{1 \leq i < j \leq n} E(w_{ij}^4).$$

We have

$$\begin{aligned}
w_{ijn}^4 &= \frac{16}{n^4} \left(\sum_{l=1}^n \frac{a_l}{\lambda_l} \operatorname{Re} \langle Y_i, \phi_l \rangle \overline{\langle Y_j, \phi_l \rangle} \right)^4 \\
&\leq \frac{16}{n^4} \left(\sum_{l=1}^n \frac{a_l}{\lambda_l} |\langle Y_i, \phi_l \rangle \langle Y_j, \phi_l \rangle| \right)^4 \\
&= \frac{16}{n^4} \sum_{l=1}^n \frac{a_l^4}{\lambda_l^4} |\langle Y_i, \phi_l \rangle|^4 |\langle Y_j, \phi_l \rangle|^4 \tag{A11a}
\end{aligned}$$

$$+ \frac{16}{n^4} \sum_{l \neq l'} \frac{a_l^3}{\lambda_l^3} \frac{a_{l'}}{\lambda_{l'}} |\langle Y_i, \phi_l \rangle|^3 |\langle Y_j, \phi_l \rangle|^3 |\langle Y_i, \phi_{l'} \rangle| |\langle Y_j, \phi_{l'} \rangle| \tag{A11b}$$

$$+ \frac{16}{n^4} \sum_{l \neq l'} \frac{a_l^2}{\lambda_l^2} \frac{a_{l'}^2}{\lambda_{l'}^2} |\langle Y_i, \phi_l \rangle|^2 |\langle Y_j, \phi_l \rangle|^2 |\langle Y_i, \phi_{l'} \rangle|^2 |\langle Y_j, \phi_{l'} \rangle|^2. \tag{A11c}$$

Consider (A11a): Using $|\langle Y_i, \phi_l \rangle|^2 \leq 4$ as before, we get $E |\langle Y_i, \phi_l \rangle|^4 \leq 4E |\langle Y_i, \phi_l \rangle|^2 = 4\lambda_l$.

Therefore,

$$E \sum_{l=1}^n \frac{a_l^4}{\lambda_l^4} |\langle Y_i, \phi_l \rangle|^4 |\langle Y_j, \phi_l \rangle|^4 \leq 16 \sum_{l=1}^n \frac{a_l^4}{\lambda_l^2}$$

and

$$\begin{aligned}
\sum_{l=1}^n \frac{a_l^4}{\lambda_l^2} &= \sum_{l=1}^n \frac{\lambda_l^6}{(\lambda_l^2 + \alpha)^4} \\
&= \sum_{l=1}^n \frac{\lambda_l^4}{(\lambda_l^2 + \alpha)^2} \frac{\lambda_l^2}{(\lambda_l^2 + \alpha)^2} \\
&\leq \sum_{l=1}^n \frac{\lambda_l^2}{(\lambda_l^2 + \alpha)^2} \leq \frac{1}{\alpha^2} \sum_{l=1}^n \lambda_l^2.
\end{aligned}$$

Hence,

$$\frac{\sum_{1 \leq i < j \leq n} (A11a)}{q_n^2} \leq \frac{C}{\alpha^2 n^2 q_n^2} \rightarrow 0$$

by the assumption $p_n / (q_n n \alpha) \rightarrow 0$.

Consider (A11b):

$$E(A11b) = \frac{16}{n^4} \sum_{l \neq l'} \frac{a_l^3}{\lambda_l^3} \frac{a_{l'}}{\lambda_{l'}} E[|\langle Y_i, \phi_l \rangle|^3 |\langle Y_i, \phi_{l'} \rangle|] E[|\langle Y_j, \phi_l \rangle|^3 |\langle Y_j, \phi_{l'} \rangle|].$$

By Cauchy-Schwarz,

$$\begin{aligned}
E[|\langle Y_i, \phi_l \rangle|^3 |\langle Y_i, \phi_{l'} \rangle|] &\leq \sqrt{E[|\langle Y_i, \phi_l \rangle|^6] E[|\langle Y_i, \phi_{l'} \rangle|^2]} \\
&\leq 4 \sqrt{E[|\langle Y_i, \phi_l \rangle|^2] E[|\langle Y_i, \phi_{l'} \rangle|^2]} \\
&= 4 \sqrt{\lambda_l} \sqrt{\lambda_{l'}}.
\end{aligned}$$

Hence,

$$E(A11b) \leq \frac{C}{n^4} \sum_{l \neq l'} \frac{a_l^3}{\lambda_l^2} a_{l'} \leq C \frac{p_n}{n^4} \sum_l \frac{a_l^3}{\lambda_l^2}.$$

Moreover,

$$\sum_l \frac{a_l^3}{\lambda_l^2} = \sum_l \frac{\lambda_l^4}{(\lambda_l^2 + \alpha)^3} \leq \sum_l \frac{1}{(\lambda_l^2 + \alpha)} \leq \frac{n}{\alpha}.$$

It follows that

$$\frac{\sum_{1 \leq i < j \leq n} (A11b)}{q_n^2} \leq C \frac{p_n}{q_n^2 n \alpha} \rightarrow 0.$$

Now, consider (A11c):

$$\begin{aligned} E(A11c) &= \frac{16}{n^4} \sum_{l \neq l'} \frac{a_l^2}{\lambda_l^2} \frac{a_{l'}^2}{\lambda_{l'}^2} E[|\langle Y_i, \phi_l \rangle|^2 |\langle Y_i, \phi_{l'} \rangle|^2] E[|\langle Y_j, \phi_l \rangle|^2 |\langle Y_j, \phi_{l'} \rangle|^2] \\ &\leq \frac{C}{n^4} \sum_{l \neq l'} \frac{a_l^2}{\lambda_l^2} \frac{a_{l'}^2}{\lambda_{l'}^2} \lambda_l \lambda_{l'} \\ &\leq \frac{C}{n^4} \left(\sum_l \frac{a_l^2}{\lambda_l} \right)^2. \end{aligned}$$

Moreover,

$$\begin{aligned} \sum_l \frac{a_l^2}{\lambda_l} &= \sum_l \frac{\lambda_l^3}{(\lambda_l^2 + \alpha)^2} = \sum_l \frac{\lambda_l^2}{(\lambda_l^2 + \alpha)} \frac{\lambda_l}{(\lambda_l^2 + \alpha)} \\ &\leq \sum_l \frac{\lambda_l}{(\lambda_l^2 + \alpha)} \leq \frac{\sum_l \lambda_l}{\alpha}. \end{aligned}$$

Therefore,

$$\frac{\sum_{1 \leq i < j \leq n} (A11c)}{q_n^2} \leq \frac{C}{\alpha^2 n^2 q_n^2} \rightarrow 0.$$

It follows that $G_I = o(\sigma(n)^4)$.

Now consider G_{II} :

$$\begin{aligned} E(w_{ijn}^2 w_{ikn}^2) &\leq \frac{16}{n^4} E \left[\left(\sum_{l=1}^n \frac{a_l}{\lambda_l} |\langle Y_i, \phi_l \rangle \langle Y_j, \phi_l \rangle| \right)^2 \left(\sum_{l'=1}^n \frac{a_{l'}}{\lambda_{l'}} |\langle Y_i, \phi_{l'} \rangle \langle Y_k, \phi_{l'} \rangle| \right)^2 \right] \\ &= \frac{16}{n^4} \sum_{l, l'} \frac{a_l^2 a_{l'}^2}{\lambda_l^2 \lambda_{l'}^2} E[|\langle Y_i, \phi_l \rangle|^2 |\langle Y_j, \phi_l \rangle|^2 |\langle Y_i, \phi_{l'} \rangle|^2 |\langle Y_k, \phi_{l'} \rangle|^2] \end{aligned}$$

because the cross products equal zero. We have

$$\begin{aligned} E[|\langle Y_i, \phi_l \rangle|^2 |\langle Y_j, \phi_l \rangle|^2 |\langle Y_i, \phi_{l'} \rangle|^2 |\langle Y_k, \phi_{l'} \rangle|^2] &\leq 4E[|\langle Y_i, \phi_l \rangle|^2] E[|\langle Y_j, \phi_l \rangle|^2] E[|\langle Y_i, \phi_{l'} \rangle|^2] E[|\langle Y_k, \phi_{l'} \rangle|^2] \\ &= 4\lambda_l^2 \lambda_{l'}^2. \end{aligned}$$

Hence,

$$E(w_{ijn}^2 w_{ikn}^2) \leq \frac{C}{n^4} \sum_l a_l^2 \sum_{l'} \frac{a_{l'}^2}{\lambda_{l'}} \leq \frac{C q_n}{n^4} \frac{\sum \lambda_l}{\alpha},$$

$$\frac{\sum_{1 \leq i < j \leq n} E[w_{ijn}^2 w_{ikn}^2]}{q_n^2} \leq \frac{C}{n^2 \alpha q_n} \rightarrow 0.$$

The other terms of G_{II} have the same form. Therefore, $G_{II} = o(\sigma(n)^4)$.

Consider G_{IV} :

$$\begin{aligned} & E(w_{ijn} w_{ikn} w_{mjn} w_{mkn}) \\ & \leq \frac{16}{n^4} E \left[\left| \sum_{l=1}^n \frac{a_l}{\lambda_l} \langle Y_i, \phi_l \rangle \overline{\langle Y_j, \phi_l \rangle} \right| \left| \sum_{l'=1}^n \frac{a_{l'}}{\lambda_{l'}} \langle Y_i, \phi_{l'} \rangle \overline{\langle Y_k, \phi_{l'} \rangle} \right| \right. \\ & \quad \left. \left| \sum_{g=1}^n \frac{a_g}{\lambda_g} \langle Y_m, \phi_g \rangle \overline{\langle Y_j, \phi_g \rangle} \right| \left| \sum_{g'=1}^n \frac{a_{g'}}{\lambda_{g'}} \langle Y_m, \phi_{g'} \rangle \overline{\langle Y_k, \phi_{g'} \rangle} \right| \right] \\ & = \frac{16}{n^4} \sum_{l=1}^n \frac{a_l^4}{\lambda_l^4} E[|\langle Y_i, \phi_l \rangle|^2 |\langle Y_j, \phi_l \rangle|^2 |\langle Y_m, \phi_l \rangle|^2 |\langle Y_k, \phi_l \rangle|^2] \end{aligned}$$

because $\langle Y_i, \phi_l \rangle$, $l = 1, 2, \dots$ are uncorrelated across l . As Y_i , $i = 1, 2, \dots$ are iid, we have

$$\begin{aligned} E(w_{ijn} w_{ikn} w_{mjn} w_{mkn}) & = \frac{16}{n^4} \sum_{l=1}^n \frac{a_l^4}{\lambda_l^4} E[|\langle Y_i, \phi_l \rangle|^2]^4 \\ & = \frac{16}{n^4} \sum_{l=1}^n a_l^4 \\ & \leq \frac{16 q_n}{n^4}. \end{aligned}$$

It follows that

$$\frac{1}{q_n^2} \sum_{1 \leq i < j < k < m \leq n} E(w_{ijn} w_{ikn} w_{mjn} w_{mkn}) \leq \frac{16}{q_n} \rightarrow 0.$$

As the other terms in G_{IV} have the same form, we can conclude that $G_{IV} = o(\sigma(n)^4)$.

Therefore, all the conditions of Lemma 16 are satisfied and the result follows. \square

Lemma 18 *We have under H_{2n} :*

$$\|K_n - \widehat{K}\| = O_{p_{2n}} \left(\max \left(n^{-1/2}, b_n^2 \right) \right), \quad (\text{A12})$$

$$\left\| (K_n^\alpha)^{-1/2} - (\widehat{K}^\alpha)^{-1/2} \right\| = O_{p_{2n}} \left(\alpha^{-3/4} \max \left(n^{-1/2}, b_n^2 \right) \right). \quad (\text{A13})$$

Proof of Lemma 18. We have $\|K_n - \widehat{K}\| \leq \|K_n - \widehat{K}\|_{HS}$ where $\|\cdot\|_{HS}$ denotes the

Hilbert-Schmidt norm. Moreover

$$\left\| K_n - \widehat{K} \right\|_{HS} \leq \|K_n - EK_n\|_{HS} + \left\| EK_n - E\widehat{K} \right\|_{HS} + \left\| E\widehat{K} - \widehat{K} \right\|_{HS}.$$

The kernel of EK_n is $\psi_n(s-t) - \psi_n(s)\psi_n(-t)$ and the kernel of $E\widehat{K}$ is

$$\begin{aligned} & E \left[(e^{isX_i} - \psi_0(s)) (e^{-itX_i} - \psi_0(-t)) \right] \\ &= \psi_n(s-t) - \psi_n(s)\psi_0(-t) - \psi_0(s)\psi_n(-t) + \psi_0(s)\psi_0(-t). \end{aligned}$$

Hence

$$\begin{aligned} EK_n - E\widehat{K} &= [\psi_n(s) - \psi_0(s)] [\psi_0(-t) - \psi_n(-t)] \\ &= -b_n^2 \eta(s) \eta(-t). \end{aligned}$$

Therefore $\left\| EK_n - E\widehat{K} \right\|_{HS} = O(b_n^2)$. Using a proof similar to that of Theorem 4 in Carrasco and Florens (2000), we have

$$\begin{aligned} \|K_n - EK_n\|_{HS} &= O_{p_{2n}} \left(\frac{1}{\sqrt{n}} \right), \\ \left\| E\widehat{K} - \widehat{K} \right\|_{HS} &= O_{p_{2n}} \left(\frac{1}{\sqrt{n}} \right). \end{aligned}$$

Hence, the result of (A12) follows. The result of (A13) can be established using a proof similar to that of Lemma B.2 in Carrasco et al (2007).

Proof of Proposition 10. As in Carrasco and Florens (2000, proof of Theorem 10), the proof proceeds in three steps.

Step 1. Let P_n denote the projection which associates to an operator K the operator K_2 defined by the first n eigenvalues and eigenfunctions of K . We show that

$$\frac{1}{\sqrt{q_n}} \left\{ \left\| (\widehat{K}^\alpha)^{-1/2} \sqrt{n} \widehat{h} \right\| - \left\| P_n (K_n^\alpha)^{-1/2} \sqrt{n} \widehat{h} \right\| \right\} \xrightarrow{P} 0 \quad (\text{A14})$$

under H_{2n} .

First note that

$$\begin{aligned} \sqrt{n} \widehat{h} &= \frac{1}{\sqrt{n}} \sum_{i=1}^n (e^{itX_i} - \psi_0(t)) \\ &= \frac{1}{\sqrt{n}} \sum_{i=1}^n (e^{itX_i} - \psi_n(t)) + \sqrt{n} (\psi_n(t) - \psi_0(t)) \\ &= O_{p_{2n}}(1) + \sqrt{n} b_n \eta(t). \end{aligned}$$

Hence $\left\| \sqrt{n}\hat{h} \right\| = O_{p_{2n}}(\sqrt{nb_n})$. We have

$$\begin{aligned}
& \frac{1}{\sqrt{q_n}} \left\{ \left\| (\hat{K}^\alpha)^{-1/2} \sqrt{n}\hat{h} \right\| - \left\| P_n (K_n^\alpha)^{-1/2} \sqrt{n}\hat{h} \right\| \right\} \\
& \leq \frac{1}{\sqrt{q_n}} \left\| [(\hat{K}^\alpha)^{-1/2} - P_n (K_n^\alpha)^{-1/2}] \sqrt{n}\hat{h} \right\| \\
& \leq \frac{1}{\sqrt{q_n}} \|P_n\| \left\| (\hat{K}^\alpha)^{-1/2} - (K_n^\alpha)^{-1/2} \right\| \left\| \sqrt{n}\hat{h} \right\| \\
& = O_{p_{2n}} \left(\frac{\max(b_n, \sqrt{nb_n^3})}{\sqrt{q_n} \alpha^{3/4}} \right)
\end{aligned}$$

because $\|P_n\| \leq 1$, $\left\| \sqrt{n}\hat{h} \right\| = O_{p_{2n}}(\sqrt{nb_n})$ and $\left\| (\hat{K}^\alpha)^{-1/2} - (K_n^\alpha)^{-1/2} \right\| = O_{p_{2n}}(\alpha^{-3/4} \max(n^{-1/2}, b_n^2))$ by Lemma 18. Therefore (A14) is satisfied.

Step 2. Show that

$$\hat{p}_n - p_{n,n} \xrightarrow{P} 0 \text{ and } \hat{q}_n - q_{n,n} \xrightarrow{P} 0$$

under H_{2n} as $n\alpha^2 \rightarrow \infty$ and $b_n^2/n \rightarrow 0$.

Using the proofs of Theorems 4 and 10 in Carrasco and Florens (2000), we can show that $\hat{p}_n - p_{n,n} = O_p\left(\frac{\|\hat{K} - K_n\|}{\alpha}\right)$ and $\hat{q}_n - q_{n,n} = O_p\left(\frac{\|\hat{K} - K_n\|}{\alpha}\right)$.

Step 3. By Lemma 17, we have under H_{2n}

$$\frac{\left\| P_n (K_n^\alpha)^{-1/2} \sqrt{n}\hat{h} \right\| - p_{n,n}}{\sqrt{q_{n,n}}} = \frac{\sum_{l=1}^n \frac{a_{l,n}}{\lambda_{l,n}} \left| \left\langle \sqrt{n}\hat{h}, \phi_{l,n} \right\rangle \right|^2 - p_{n,n}}{\sqrt{q_{n,n}}} \xrightarrow{d} \mathcal{N}(c \|\eta\|_K^2, 1).$$

Using steps 1 and 2, we obtain the desired result. \square

Proof of Proposition 11. Let $\{\phi_j, \lambda_j\}$ be the eigenfunctions and eigenvalues of K_θ . Let ψ_j such that $\phi_j = U_\theta^* \psi_j$ and consequently $U_\theta \phi_j = U_\theta U_\theta^* \psi_j = \psi_j$. We have

$$\begin{aligned}
U_\theta K_\theta U_\theta^* \psi_j &= U_\theta K_\theta \phi_j \\
&= \lambda_j U_\theta \phi_j \\
&= \lambda_j \psi_j.
\end{aligned}$$

Therefore, $\{\psi_j, \lambda_j\}$ are the eigenfunctions and eigenvalues of \tilde{K} . It follows that

$$\begin{aligned}
\|U_\theta \hat{h}(\theta)\|_{(U_\theta K_\theta U_\theta^*)^\alpha}^2 &= \sum_j \frac{\lambda_j}{\lambda_j^2 + \alpha} \left| \langle U_\theta \hat{h}(\theta), \psi_j \rangle \right|^2 \\
&= \sum_j \frac{\lambda_j}{\lambda_j^2 + \alpha} \left| \langle \hat{h}(\theta), U_\theta^* \psi_j \rangle \right|^2 \\
&= \sum_j \frac{\lambda_j}{\lambda_j^2 + \alpha} \left| \langle \hat{h}(\theta), \phi_j \rangle \right|^2 \\
&= \|\hat{h}(\theta)\|_{K_\theta^\alpha}^2,
\end{aligned}$$

as desired. □

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Table 1: Bootstrap-based size corrected rejection rates at 1%, 5%, and 10% significance levels for univariate Gaussian null hypothesis

ω	$J(\hat{\theta}, K_{\hat{\theta}})$						$J(\hat{\theta}, \hat{K})$						LM tests							
	$\alpha = .1$		$\alpha = .01$		$\alpha = 0.01$		$\alpha = 0.01$		$\alpha = 0.01$		$\alpha = 0.01$									
	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1	KS	CvM	AD	BW	S	t	A	t
Panel A: Student t with 12 degrees of freedom																				
10%	15.0	21.9	15.0	23.9	26.8	8.8	28.6	8.5	15.1	16.9	20.1	12.5	31.4	32.3						
5%	8.0	13.9	8.0	15.5	16.8	4.3	19.3	4.7	8.8	9.5	11.7	6.4	25.4	24.3						
1%	2.5	4.5	2.3	5.9	5.5	1.1	7.0	1.3	1.7	3.1	4.1	1.1	12.9	12.4						
Panel B: Scale mixture of two normals. the outliers case																				
10%	16.3	23.4	15.7	24.2	23.9	11.9	26.0	11.6	17.5	18.7	20.3	13.3	30.9	33.3						
5%	9.4	15.3	9.1	15.6	14.8	6.4	17.5	6.3	10.8	11.3	12.3	7.1	24.8	25.7						
1%	3.0	5.5	2.8	6.4	5.2	1.7	6.3	1.7	2.9	4.0	4.8	1.6	13.6	14.0						
Panel C: Scale mixture of two normals, the inliers case																				
10%	21.2	21.2	19.8	22.9	23.8	5.1	19.2	4.0	21.4	21.7	21.4	7.8	21.7	23.7						
5%	13.4	13.3	12.6	14.6	14.8	2.4	12.1	1.6	14.0	13.9	13.5	3.5	15.6	15.4						
1%	5.4	4.1	4.7	5.3	4.6	0.4	3.4	0.2	4.2	5.3	5.3	0.7	5.5	4.9						

Notes: Results based on 10,000 samples of size $n = 100$. Critical values are computed using parametric bootstrap. $J(\hat{\theta}, K_{\hat{\theta}})$ and $J(\hat{\theta}, \hat{K})$ denote the test that uses the theoretical covariance operator and the centred sample estimator of the covariance operator, respectively. α denotes the regularization parameter and ω is the scale parameter of the $\mathcal{N}(0, \omega^2)$ density defining inner products. S t and A t are the Lagrange multiplier tests based on the score of the symmetric and asymmetric Student t , respectively. KS, CvM and AD denote the Kolmogorov–Smirnov, the Cramér–von Mises and the Anderson–Darling tests while BW denotes the test in Bierens and Wang (2012).

Table 1 (cont.): Bootstrap-based size corrected rejection rates at 1%, 5%, and 10% significance levels for univariate Gaussian null hypothesis

ω	$J(\hat{\theta}, K_{\hat{\theta}})$										$J(\hat{\theta}, \hat{K})$			LM tests			
	$\alpha = .1$			$\alpha = .01$			$\alpha = 0.01$			Other consistent tests			S t A t				
	$\sqrt{10}$	1	1	$\sqrt{10}$	1	1	$\sqrt{10}$	1	1	KS	CvM	AD	BW	S	t	A	t
Panel D: Asymmetric Student t with 12 degrees of freedom and $\beta = -0.75$																	
10%	25.6	43.8	24.2	43.2	35.3	13.9	39.1	17.1	28.7	33.7	37.7	17.4	40.9	51.3			
5%	16.6	34.1	15.3	33.1	24.5	8.0	29.4	11.2	19.7	23.3	27.5	10.3	34.9	41.6			
1%	6.2	17.4	5.3	17.4	9.6	2.8	13.3	4.9	6.7	10.8	13.9	3.0	21.8	25.2			
Panel E: Asymmetric (but with zero skewness), mesokurtic, location-scale mixture of three normals																	
10%	16.7	25.3	20.4	25.1	13.4	29.4	20.9	37.3	15.6	20.3	25.9	12.6	21.1	22.5			
5%	7.7	14.7	11.4	13.9	7.7	15.6	13.0	23.4	7.6	9.9	14.4	6.5	9.0	9.7			
1%	1.4	3.6	2.8	2.8	2.6	3.7	3.7	6.5	1.1	1.8	3.7	1.2	1.1	2.4			
Panel F: Standardized second order Hermite expansion of the standard normal																	
10%	45.2	57.7	40.6	56.0	16.4	18.2	20.7	31.0	45.1	52.0	53.3	24.7	9.2	38.0			
5%	32.3	45.3	28.7	42.6	9.0	10.6	12.8	19.5	32.8	39.1	40.2	16.1	3.2	18.6			
1%	14.9	20.8	11.9	19.3	2.2	3.0	3.6	6.7	12.5	19.4	20.1	5.2	0.4	2.2			

Results based on 10,000 samples of size $n = 100$. Critical values are computed using parametric bootstrap. $J(\hat{\theta}, K_{\hat{\theta}})$ and $J(\hat{\theta}, \hat{K})$ denote the test that uses the theoretical covariance operator and the centred sample estimator of the covariance operator, respectively. α denotes the regularization parameter and ω is the scale parameter of the $\mathcal{N}(0, \omega^2)$ density defining inner products. S t and A t are the Lagrange multiplier tests based on the score of the symmetric and asymmetric Student t , respectively. KS, CvM and AD denote the Kolmogorov–Smirnov, the Cramér–von Mises and the Anderson–Darling tests while BW denotes the test in Bierens and Wang (2012).

Table 2: Bootstrap-based size corrected rejection rates at 1%, 5%, and 10% significance levels for univariate Uniform null hypothesis

ω	$J(\theta_0, K)$						$J(\theta_0, \hat{K})$						Directional tests					
	$\alpha = .1$		$\alpha = .01$		$\alpha = 0.01$		$\alpha = 0.01$		$\alpha = 0.01$		$\alpha = 0.01$							
	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1	KS	CvM	AD	BW	LM	BM
Panel A: Symmetric Beta with parameters $a = b = 1.1$																		
10%	22.6	24.6	25.1	25.5	17.7	15.0	21.5	20.9	10.4	10.3	10.0	10.7	12.6	13.7				
5%	13.5	14.6	15.2	15.9	9.5	7.6	11.9	11.8	5.0	4.5	4.4	5.5	5.6	7.7				
1%	4.0	4.5	4.5	5.0	2.1	1.7	3.1	3.1	0.9	0.8	0.8	1.0	0.8	1.6				
Panel B: Asymmetric Beta with parameters $a = 1.1$ and $b = 1$																		
10%	19.3	19.2	19.5	17.3	17.4	18.2	19.4	21.9	18.9	20.9	20.7	15.0	18.2	18.3				
5%	11.5	11.0	11.4	9.8	9.7	10.9	11.1	13.2	11.7	12.4	12.4	8.5	9.8	11.2				
1%	3.4	3.3	3.1	2.5	2.3	3.3	2.9	4.1	2.9	3.4	3.7	2.6	2.5	2.8				
Panel C: Standard Gaussian PITs of obs. drawn from a univariate asymmetric Student t with 12 df and $\beta = -.75$																		
10%	27.4	26.8	26.4	27.1	16.6	11.6	18.7	18.5	13.8	13.7	14.5	11.8	16.6	14.2				
5%	17.5	16.3	16.9	16.7	9.4	6.5	10.8	10.9	8.0	6.8	6.8	6.4	9.6	8.1				
1%	6.6	5.8	6.1	5.9	2.6	1.7	3.3	3.4	1.6	1.4	1.3	1.2	3.8	1.7				

Results based on 10,000 samples of size $n = 100$. Critical values are computed using parametric bootstrap. $J(\theta_0, K)$ and $J(\theta_0, \hat{K})$ denote the test that uses the theoretical covariance operator and the centred sample estimator of the covariance operator, respectively. α denotes the regularization parameter and ω is the scale parameter of the $\mathcal{N}(0, \omega^2)$ density defining inner products. LM is the LM test of uniform vs beta proposed by Selton (1992) and BM is a moment test based on the first two Jacobi polynomials proposed by Bontemps and Meddahi (2012). KS, CvM and AD denote the Kolmogorov–Smirnov, the Cramér–von Mises and the Anderson–Darling tests while BW denotes the test in Bierens and Wang (2012).

Table 3: Bootstrap-based size corrected rejection rates at 1%, 5%, and 10% significance levels for bivariate Gaussian null hypothesis

ω	$J(\hat{\theta}, K_{\hat{\theta}})$			$J(\hat{\theta}, \hat{K})$			LM tests	
	$\alpha = .1$	$\alpha = .01$	$\alpha = 0.01$	$\alpha = 0.01$	$\alpha = 0.01$	$\alpha = 0.01$	S	t
	$\sqrt{10}$	1	$\omega = 1$	$\sqrt{10}$	1	$\sqrt{10}$	A	t
Panel A: Student t with 12 degrees of freedom								
10%	20.5	30.1	20.0	32.1	31.0	17.1	35.5	43.7
5%	13.1	20.3	13.1	22.7	19.7	10.1	23.7	33.4
1%	3.8	8.5	3.7	11.1	5.4	3.4	8.0	16.8
Panel B: Scale mixture of two normals, the outliers case								
10%	16.8	27.2	16.5	28.4	23.7	22.1	27.4	44.5
5%	10.1	18.1	9.9	20.0	14.2	13.6	17.1	33.7
1%	2.5	7.2	2.4	9.3	3.5	4.6	5.1	17.2
Panel C: Asymmetric Student t with 12 degrees of freedom and $\beta_1 = \beta_2 = -0.75$								
10%	39.7	65.9	38.4	64.7	40.8	33.1	46.7	63.6
5%	29.1	55.6	28.7	54.9	27.7	23.5	34.0	54.0
1%	12.1	37.2	11.6	38.8	9.2	11.3	14.0	36.1

Results based on 10,000 samples of size $n = 100$. Critical values are computed using parametric bootstrap. $J(\hat{\theta}, K_{\hat{\theta}})$ and $J(\hat{\theta}, \hat{K})$ denote the test that uses the theoretical covariance operator and the centred sample estimator of the covariance operator, respectively. α denotes the regularization parameter and ω is the scale parameter of the $\mathcal{N}(0, \omega^2)$ density defining inner products. S t and A t are the Lagrange multiplier tests based on the score of the symmetric and asymmetric Student t , respectively.

Table 4: Bootstrap-based size corrected rejection rates at 1%, 5%, and 10% significance levels for univariate χ^2_2 null hypothesis

ω	$J(\theta_0, K)$						$J(\theta_0, \hat{K})$						LM tests	
	$\alpha = .1$		$\alpha = .01$		$\alpha = .1$		$\alpha = .01$		$\alpha = .1$		$\alpha = .01$		F	G
	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1		
Panel A: Scaled F distribution with 2 and 12 degrees of freedom														
10%	31.9	35.6	31.7	28.5	20.7	7.3	13.0	5.4	22.8	23.4	25.1	14.9	51.8	34.1
5%	20.0	22.8	20.3	17.4	13.5	4.7	6.9	3.2	14.3	14.4	15.4	8.5	43.9	25.6
1%	6.1	7.8	6.2	5.7	4.6	1.6	2.2	1.1	4.3	4.4	4.7	2.3	26.5	12.0
Panel B: Gamma distribution with shape parameter $\alpha = 2/3$ and scale parameter $\beta = 3$														
10%	29.0	41.8	31.8	27.0	22.1	5.6	7.2	3.5	59.3	61.3	80.0	29.6	67.4	95.7
5%	17.9	29.0	20.7	16.8	15.3	4.0	3.7	1.9	45.7	46.2	68.1	18.9	56.2	93.2
1%	5.6	12.0	6.6	5.7	6.4	1.4	1.3	0.8	22.1	22.2	41.8	6.0	28.8	84.0
Panel C: Square norm of observations drawn from a bivariate asymmetric Student t with 12 df and $\beta = -.75\ell$														
10%	40.6	47.1	40.5	39.4	28.7	9.2	18.7	6.9	29.0	29.8	32.1	17.9	62.1	40.9
5%	27.1	32.5	27.6	26.1	19.4	7.0	11.4	4.6	18.7	20.0	21.3	10.4	55.6	32.3
1%	9.3	13.2	9.4	9.9	7.3	3.2	4.1	2.2	6.5	7.1	7.9	3.3	39.7	17.9

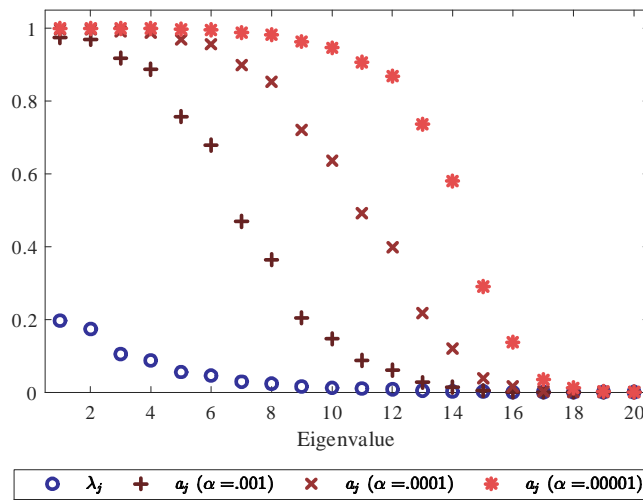
Results based on 10,000 samples of size $n = 100$. Critical values are computed using parametric bootstrap. $J(\theta_0, K)$ and $J(\theta_0, \hat{K})$ denote the test that uses the theoretical covariance operator and the centred sample estimator of the covariance operator, respectively. α denotes the regularization parameter and ω is the scale parameter of the $U[-\omega, \omega]$ density defining inner products. F is the LM test of chi square with N degrees of freedom versus F with the same number of degrees of freedom in the numerator but degrees of freedom in the denominator proposed by Fiorentini, Sentana and Calzolari (2003) and G the LM test against a gamma distribution with mean N but shape parameter $\alpha \neq N/2$ developed in Amengual and Sentana (2012). KS, CvM and AD denote the Kolmogorov–Smirnov, the Cramér–von Mises and the Anderson–Darling tests while BW denotes the test in Bierens and Wang (2012).

Table 5: Bootstrap-based size corrected rejection rates at 1%, 5%, and 10% significance levels for Cauchy null hypothesis

ω	$J(\hat{\theta}, K_{\hat{\theta}})$			$J(\hat{\theta}, \hat{K})$			Other consistent tests					
	$\alpha = .1$		$\alpha = .01$	$\alpha = 0.1$		$\alpha = 0.01$	KS	CvM	AD	BW		
	$\sqrt{10}$	1	$\sqrt{10}$	1	$\sqrt{10}$	1						
Panel A: Student t with 2 degrees of freedom												
10%	100.00	42.13	100.00	53.87	89.68	40.15	90.92	57.39	20.30	12.11	43.95	55.51
5%	100.00	30.74	100.00	41.29	81.81	29.43	83.08	43.89	9.45	4.29	21.57	40.95
1%	100.00	11.95	100.00	19.49	56.89	10.70	58.95	21.05	1.30	0.44	2.29	16.52
Panel B: : Asymmetric Student t with 6 degrees of freedom and $\beta = -0.25$												
10%	100.00	77.32	100.00	92.54	99.98	74.80	99.99	92.06	99.83	45.54	100.00	93.07
5%	100.00	67.30	100.00	87.29	99.94	64.00	99.94	86.24	97.35	19.92	100.00	86.91
1%	100.00	40.93	100.00	69.85	99.48	39.76	99.55	68.06	59.92	1.99	98.40	64.66
Panel C: Laplace with location and scale parameters 0 and $1/\sqrt{2}$, respectively												
10%	100.00	27.01	100.00	44.78	96.52	21.34	97.31	43.29	99.67	10.12	100.00	74.13
5%	100.00	17.62	100.00	33.78	92.85	13.83	93.98	31.12	97.27	3.34	100.00	59.56
1%	100.00	5.28	100.00	15.42	77.84	4.09	80.43	13.49	67.67	0.32	99.12	30.87

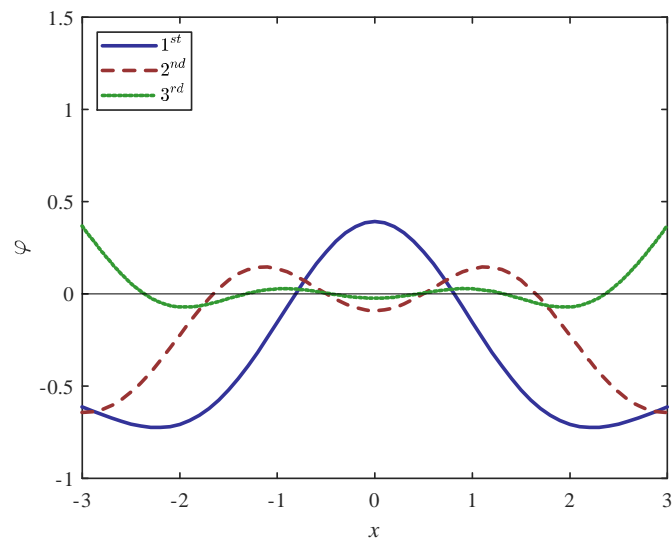
Notes: Results based on 10,000 samples of size $n = 100$. Critical values are computed using parametric bootstrap. $J(\hat{\theta}, K_{\hat{\theta}})$ and $J(\hat{\theta}, \hat{K})$ denote the test that uses the theoretical covariance operator and the centred sample estimator of the covariance operator, respectively. α denotes the regularization parameter and ω is the scale parameter of the $\mathcal{N}(0, \omega^2)$ density defining inner products. KS, CvM and AD denote the Kolmogorov-Smirnov, the Cramér-von Mises and the Anderson-Darling tests while BW denotes the test in Biersens and Wang (2012).

Figure 1: Eigenvalues (λ_j 's) and weights (a_j 's) of the covariance K for the standard Normal distribution



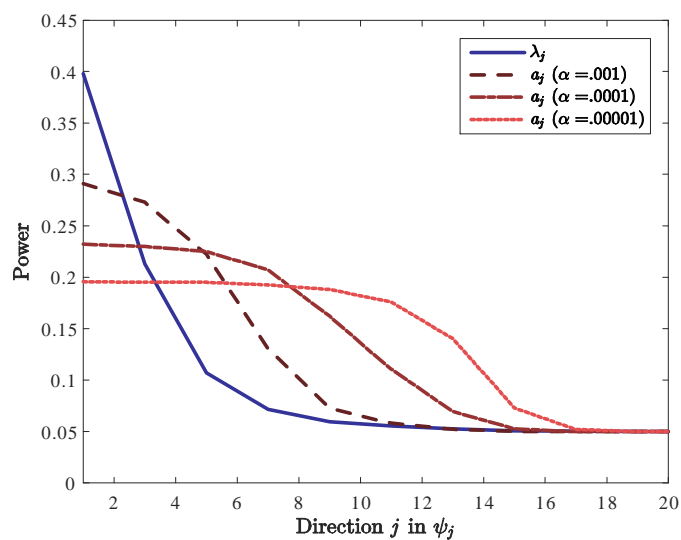
Notes: Eigenvalues are computed following the procedure described in Appendix A.1 with a grid of 1,000 points.

Figure 2: Alternatives with maximum power (φ_j , for $j = 1, 2, 3$) for the standard Normal distribution



Notes: Eigenvalues are computed following the procedure described in Appendix A.1 with a grid of 1,000 points.

Figure 3: Asymptotic power at the 5% level of the T_B tests based on $B = I$ (λ_j 's) and $B = K_\alpha^{-1/2}$ (a_j 's) under local alternatives



Notes: Eigenvalues are computed following the procedure described in Appendix A.1 with a grid of 1,000 points. Power is computed using rejection rates obtained from simulated samples of size 100,000 under both the null and the alternatives.