# BAYESIAN SEMI-PARAMETRIC ESTIMATION OF THE LONG-MEMORY PARAMETER UNDER FEXP-PRIORS. * 

By Willem Kruijer<br>Biometris, Wageningen University and Research Centre<br>AND

By Judith Rousseau
Université Paris Dauphine, and CREST-ENSAE, Paris
For a Gaussian time series with long-memory behavior, we use the FEXP-model for semi-parametric estimation of the long-memory parameter $d$. The true spectral density $f_{o}$ is assumed to have longmemory parameter $d_{o}$ and a FEXP-expansion of Sobolev-regularity $\beta>1$. We prove that when $k$ follows a Poisson or geometric prior, or a sieve prior increasing at rate $n^{\frac{1}{1+2 \beta}}$, $d$ converges to $d_{o}$ at a suboptimal rate. When the sieve prior increases at rate $n^{\frac{1}{2 \beta}}$ however, the minimax rate is almost obtained. Our results can be seen as a Bayesian equivalent of the result which Moulines and Soulier obtained for some frequentist estimators.

1. Introduction. Let $X_{t}, t \in \mathbb{Z}$, be a stationary Gaussian time series with zero mean and spectral density $f_{o}(x), x \in[-\pi, \pi]$, which takes the form

$$
\begin{equation*}
\left|1-e^{i x}\right|^{-2 d_{o}} M_{o}(x), \quad x \in[-\pi, \pi], \tag{1.1}
\end{equation*}
$$

where $d_{o} \in\left(-\frac{1}{2}, \frac{1}{2}\right)$ is called the long-memory parameter, and $M$ is a slowlyvarying bounded function that describes the short-memory behavior of the series. If $d_{o}$ is positive, this makes the autocorrelation function $\rho(h)$ decay polynomially, at rate $h^{-\left(1-2 d_{o}\right)}$, and the time series is said to have longmemory. When $d_{o}=0, X_{t}$ has short memory, and the case $d_{o}<0$ is referred to as intermediate memory. Long memory time series models are used in a wide range of applications, such as hydrological or financial time series; see for example Beran (1994) or Robinson (1994). In parametric approaches, a

[^0]finite dimensional model is used for the short memory part $M_{o}$; the most well known example is the $\operatorname{ARFIMA}(\mathrm{p}, \mathrm{d}, \mathrm{q})$ model. The asymptotic properties of maximum likelihood estimators (Dahlhaus (1989) or Lieberman et al. (2003)) and Bayesian estimators (Philippe and Rousseau (2002)) have been established in such models and these estimators are consistent and asymptotically normal with a convergence rate of order $\sqrt{n}$. However when the model for the short memory part is misspecified, the estimator for $d$ can be inconsistent, calling for semi-parametric methods for the estimation of $d$. A key feature of semi-parametric estimators of the long-memory parameter is that they converge at a rate which depends on the smoothness of the shortmemory part, and apart from the case where $M_{o}$ is infinitely smooth, the convergence rate is smaller than $\sqrt{n}$. The estimation of the long-memory parameter $d$ can thus be considered as a non-regular semi-parametric problem. In Moulines and Soulier (2003) (p. 274) it is shown that when $f_{o}$ satisfies (1.4), the minimax rate for $d$ is $n^{-\frac{2 \beta-1}{4 \beta}}$. There are frequentist estimators for $d$ based on the periodogram that achieve this rate (see Hurvich et al. (2002) and Moulines and Soulier (2003)).

Although Bayesian methods in long-memory models have been widely used (see for instance Ko et al. (2009), Jensen (2004) or Holan and McElroy (2010)), the literature on convergence properties of non- and semiparametric estimators is sparse. Rousseau et al. (2010) (RCL hereafter) obtain consistency and rates for the $L_{2}$-norm of the log-spectral densities (Theorems 3.1 and 3.2 ), but for $d$ they only show consistency (Corollary 1 ). No results exist on the posterior concentration rate on $d$, and thus on the convergence rates of Bayesian semi-parametric estimators of $d$. In this paper we aim to fill this gap for a specific family of semi-parametric priors.

We study Bayesian estimation of $d$ within the FEXP-model (Beran (1993), Robinson (1995)), that contains densities of the form

$$
\begin{equation*}
f_{d, k, \theta}(x)=\left|1-e^{i x}\right|^{-2 d} \exp \left\{\sum_{j=0}^{k} \theta_{j} \cos (j x)\right\} \tag{1.2}
\end{equation*}
$$

where $d \in\left(-\frac{1}{2}, \frac{1}{2}\right), k$ is a nonnegative integer and $\theta \in \mathbb{R}^{k+1}$. The factor $\exp \left\{\sum_{j=0}^{k} \theta_{j} \cos (j x)\right\}$ models the function $M_{o}$ in (1.1). In contrast to the original finite-dimensional FEXP-model (Beran (1993)), where $k$ was supposed to be known, or at least bounded, $f_{o}$ may have an infinite FEXPexpansion, and we allow $k$ to increase with the number of observations to obtain approximations $f$ that are increasingly close to $f_{o}$. Note that the case where the true spectral density satisfies $f_{o}=f_{d_{o}, k_{o}, \theta_{o}}$, is considered in Holan and McElroy (2010). In this paper we will pursue a fully Bayesian semi-
parametric estimation of $d$, the short memory parameter being considered as an infinite-dimensional nuisance parameter. We obtain results on the convergence rate and asymptotic distribution of the posterior distribution for $d$, which we summarize below in section 1.2. These are to our knowledge the first of this kind in the Bayesian literature on semi-parametric time series. First we state the most important assumptions.
1.1. Asymptotic framework. For observations $X=\left(X_{1}, \ldots, X_{n}\right)$ from a Gaussian stationary time series with spectral density $f$, let $T_{n}(f)$ denote the associated covariance matrix and $l_{n}(f)$ denote the log-likelihood

$$
l_{n}(f)=-\frac{n}{2} \log (2 \pi)-\frac{1}{2} \log \operatorname{det}\left(T_{n}(f)\right)-\frac{1}{2} X^{t} T_{n}^{-1}(f) X .
$$

We consider semi-parametric priors on $f$ based on the FEXP-model defined by (1.2), inducing a parametrization of $f$ in terms of $(d, k, \theta)$. Assuming priors $\pi_{d}$ for $d$, and, independent of $d, \pi_{k}$ for $k$ and $\pi_{\theta \mid k}$ for $\theta \mid k$, we study the (marginal) posterior for $d$, given by

$$
\begin{equation*}
\Pi(d \in D \mid X)=\frac{\sum_{k=0}^{\infty} \pi_{k}(k) \int_{D} \int_{\mathbb{R}^{k+1}} e^{l_{n}(d, k, \theta)} d \pi_{\theta \mid k}(\theta) d \pi_{d}(d)}{\sum_{k=0}^{\infty} \pi_{k}(k) \int_{-\frac{1}{2}}^{\frac{1}{2}} \int_{\mathbb{R}^{k+1}} e^{l_{n}(d, k, \theta)} d \pi_{\theta \mid k}(\theta) d \pi_{d}(d)} \tag{1.3}
\end{equation*}
$$

The posterior mean or median can be taken as point-estimates for $d$, but we will focuss on the posterior $\Pi(d \mid X)$ itself.

It is assumed that the true spectral density is of the form

$$
\begin{align*}
f_{o}(x) & =\left|1-e^{i x}\right|^{-2 d_{o}} \exp \left\{\sum_{j=0}^{\infty} \theta_{o, j} \cos (j x)\right\},  \tag{1.4}\\
\theta_{o} & \in \Theta\left(\beta, L_{o}\right)=\left\{\theta \in l_{2}(\mathbb{N}): \sum_{j=0}^{\infty} \theta_{j}^{2}(1+j)^{2 \beta} \leq L_{o}\right\},
\end{align*}
$$

for some known $\beta>1$.
In particular, we derive bounds on the rate at which $\Pi(d \in D \mid X)$ concentrates at $d_{o}$, together with a Bernstein - von -Mises (BVM) property of this distribution. The posterior concentration rate for $d$ is defined as the fastest sequence $\alpha_{n}$ converging to zero such that

$$
\begin{equation*}
\Pi\left(\left|d-d_{o}\right|<K \alpha_{n} \mid X\right) \xrightarrow{P_{o}} 0, \quad \text { for a given fixed } \mathrm{K} . \tag{1.5}
\end{equation*}
$$

1.2. Summary of the results. Under the above assumptions we obtain several results for the asymptotic distribution of $\Pi(d \in D \mid X)$. Our first main result (Theorem 2.1) states that under the sieve prior $k_{n} \sim(n / \log n)^{1 /(2 \beta)}$, $\Pi(d \in D \mid X)$ is asymptotically Gaussian, and we give expressions for the posterior mean and the posterior variance. A consequence (Corollary 2.1) of this result is that the convergence rate for $d$ under this prior is at least $\delta_{n}=(n / \log n)^{-\frac{2 \beta-1}{4 \beta}}$, i.e. in (1.5) $\alpha_{n}$ is bounded by $\delta_{n}$. Up to a $\log n$ term, this is the minimax rate.

By our second main result (Theorem 2.2), the rate for $d$ is suboptimal when $k$ is given a a Poisson or a Geometric distribution, or a sieve prior $k_{n}^{\prime} \sim(n / \log n)^{\frac{1}{1+2 \beta}}$. More precisely, there exists $f_{o}$ such that the posterior concentration rate $\alpha_{n}$ is greater than $n^{-(\beta-1 / 2) /(2 \beta+1)}$, and thus suboptimal. Consequently, despite having good frequentist properties for the estimation of the spectral density $f$ itself (see RCL), these priors are much less suitable for the estimation of $d$. This is not a unique phenomenon in (Bayesian) semiparametric estimation and is encountered for instance in the estimation of a linear functional of the signal in white-noise models, see Li and Zhao (2002) or Arbel (2010).

The BVM property means that asymptotically the posterior distribution of $d$ behaves like $\alpha_{n}^{-1}(d-\hat{d}) \sim \mathcal{N}(0,1)$, where $\hat{d}$ is an estimate whose frequentist distribution (associated to the parameter $d$ ) is $\mathcal{N}\left(d_{o}, \alpha_{n}^{2}\right)$. We prove such a property on the posterior distribution of $d$ given $k=k_{n}$. In regular parametric long-memory models, the BVM property has been established by Philippe and Rousseau (2002). It is however much more difficult to establish BVM theorems in infinite dimensional setups, even for independent and identically distributed models; see for instance Freedman (1999), Castillo (2010) and Rivoirard and Rousseau (2010). In particular it has been proved that the BVM property may not be valid, even for reasonable priors. The BVM property is however very useful since it induces a strong connection between frequentist and Bayesian methods. In particular, it implies that Bayesian credible regions are asymptotically also frequentist confidence regions with the same nominal level. In section 2 we discuss this issue in more detail.
1.3. Overview of the paper. In section 2, we present three families of priors based on the sieve model defined by (1.2) with either $k$ increasing at the rate $(n / \log n)^{1 /(2 \beta)}, k$ increasing at the rate $(n / \log n)^{1 /(2 \beta+1)}$ or with random $k$. We study the behavior of the posterior distribution of $d$ in each case and prove that the former leads to optimal frequentist procedures while the latter two lead to suboptimal procedures. In section 3 we give a decom-
position of $\Pi(d \in D \mid X)$ defined in (1.3), and obtain bounds for the terms in this decomposition in sections 3.2 and 3.3. Using these results we prove Theorems 2.1 and 2.2 in respectively sections 4 and 5 . Conclusions are given in section 6 . In the appendices we give the proofs of the lemmas in section 3 , as well as some additional results on the derivatives of the log-likelihood. The proofs of various technical results can be found in the supplementary material. We conclude this introduction with an overview of the notation.
1.4. Notation. The $m$-dimensional identity matrix is denoted $I_{m}$. We write $|A|$ for the Frobenius or Hilbert-Schmidt norm of a matrix $A$, i.e. $|A|=$ $\sqrt{\operatorname{tr} A A^{t}}$, where $A^{t}$ denotes the transpose of $A$. The operator or spectral norm is denoted $\|A\|^{2}=\sup _{\|x\|=1} x^{t} A^{t} A x$. We also use $\|\cdot\|$ for the Euclidean norm on $\mathbb{R}^{k}$ or $l^{2}(\mathbb{N})$. The inner-product is denoted $|\cdot|$. We make frequent use of the relations

$$
\begin{align*}
& |A B|=|B A| \leq\|A\| \cdot|B|, \quad\|A B\| \leq\|A\| \cdot\|B\|, \quad\|A\| \leq|A| \leq \sqrt{n}\|A\|,  \tag{1.6}\\
& |\operatorname{tr}(A B)|=|\operatorname{tr}(B A)| \leq|A| \cdot|B|, \quad\left|x^{t} A x\right| \leq x^{t} x\|A\|,
\end{align*}
$$

see Dahlhaus (1989), p. 1754. For any function $h \in L_{1}([-\pi, \pi]), T_{n}(h)$ is the matrix with entries $\int_{-\pi}^{\pi} e^{i l l-m \mid x} h(x) d x, l, m=1, \ldots, n$. For example, $T_{n}(f)$ is the covariance matrix of observations $X=\left(X_{1}, \ldots, X_{n}\right)$ from a time series with spectral density $f$. If $h$ is square integrable on $[-\pi, \pi]$ we note

$$
\|h\|_{2}=\int_{-\pi}^{\pi} h^{2}(x) d x .
$$

The norm $l$ between spectral densities $f$ and $g$ is defined as

$$
l(f, g)=\frac{1}{2 \pi} \int_{-\pi}^{\pi}(\log f(x)-\log g(x))^{2} d x
$$

Unless stated otherwise, all expectations and probabilities are with respect to $P_{o}$, the law associated with the true spectral density $f_{o}$. To avoid ambiguous notation (e.g. $\theta_{0}$ versus $\theta_{0,0}$ ) we write $\theta_{o}$ instead of $\theta_{0}$. Related quantities such as $f_{o}$ and $d_{o}$ are also denoted with the $o$-subscript.

The symbols $o_{P}$ and $O_{P}$ have their usual meaning. We use boldface when they are uniform over a certain parameter range. Given a probability law $P$, a family of random variables $\left\{W_{d}\right\}_{d \in A}$ and a positive sequence $a_{n}, W_{d}=$ $\mathbf{o}_{\mathbf{P}}\left(a_{n}, A\right)$ means that

$$
P\left(\sup _{d \in A}\left|W_{d}\right| / a_{n}>\epsilon\right) \rightarrow 0,(n \rightarrow \infty) .
$$

When the parameter set is clear from the context we simply write $\mathbf{o}_{\mathbf{P}}\left(a_{n}\right)$. In a similar fashion, we write $\mathbf{o}\left(a_{n}\right)$ when the sequence is deterministic. In conjunction with the $o_{P}$ and $O_{P}$ notation we use the letters $\delta$ and $\epsilon$ as follows. When, for some $\tau>0$ and a probability $P$ we write $Z=O_{P}\left(n^{\tau-\epsilon}\right)$, this means that $Z=O\left(n^{\tau+\epsilon}\right)$ for all $\epsilon>0$. When, on the other hand, $Z=O_{P}\left(n^{\tau-\delta}\right)$, we mean that this is true for some $\delta>0$. If the value of $\delta$ is of importance it is given a name, for example $\delta_{1}$ in Lemma 3.4.

The true spectral density of the process is denoted $f_{o}$. We denote $k$ dimensional Sobolev-balls by

$$
\begin{equation*}
\Theta_{k}(\beta, L)=\left\{\theta \in \mathbb{R}^{k+1}: \sum_{j=0}^{k} \theta_{j}^{2}(1+j)^{2 \beta} \leq L\right\} \subset \mathbb{R}^{k+1} \tag{1.7}
\end{equation*}
$$

For any real number $x$, let $x_{+}$denote $\max (0, x)$. The number $r_{k}$ denotes the $\operatorname{sum} \sum_{j \geq k+1} j^{-2}$. Let $\eta$ be the sequence defined by $\eta_{j}=-2 / j, j \geq 1$ and $\eta_{0}=0$. For an infinite sequence $u=\left(u_{j}\right)_{j \geq 0}$, let $u_{[k]}$ denote the vector of the first $k+1$ elements. In particular, $\eta_{[k]}=\left(\eta_{0}, \ldots, \eta_{k}\right)$. The letter $C$ denotes any generic constant independent of $L_{o}$ and $L$, which are the constants appearing in the assumptions on $f_{o}$ and the definition of the prior.
2. Main results. Before stating Theorems 2.1 and 2.2 in section 2.3, we state the assumptions on $f_{o}$ and the prior, and give examples of priors satisfying these assumptions.
2.1. Assumptions on the prior and the true spectral density. We assume observations $X=\left(X_{1}, \ldots, X_{n}\right)$ from a stationary Gaussian time series with law $P_{o}$, which is a zero mean Gaussian distribution, whose covariance structure is defined by a spectral density $f_{o}$ satisfying (1.4), for known $\beta>1$. It is assumed that for a small constant $t>0, d_{o} \in\left[-\frac{1}{2}+t, \frac{1}{2}-t\right]$.

Assumptions on $\Pi$. We consider different priors, and first state the assumptions that are common to all these priors. The prior on the space of spectral densities consists of independent priors $\pi_{d}, \pi_{k}$ and, conditional on $k$, $\pi_{\theta \mid k}$. The prior for $d$ has density $\pi_{d}$ which is strictly positive on $\left[-\frac{1}{2}+t, \frac{1}{2}-t\right]$, the interval which is assumed to contain $d_{o}$, and zero elsewhere. The prior for $\theta$ given $k$ has a density $\pi_{\theta \mid k}$ with respect to Lebesgue measure. This density satisfies condition $\operatorname{Hyp}\left(\mathcal{K}, \mathrm{c}_{0}, \beta, \mathrm{~L}_{\mathrm{o}}\right)$, by which we mean that for a subset $\mathcal{K}$ of $\mathbb{N}$,

$$
\min _{k \in \mathcal{K}} \inf _{\theta \in \Theta_{k}\left(\beta, L_{o}\right)} e^{c_{0} k \log k} \pi_{\theta \mid k}(\theta)>1
$$

where $L_{o}$ is as in (1.4). The choice of $\mathcal{K}$ depends on the prior for $k$ and $\theta \mid k$. We consider the following classes of priors.

- Prior A: $k$ is deterministic and increasing at rate

$$
\begin{equation*}
k_{n}=\left\lfloor k_{A}(n / \log n)^{\frac{1}{2 \beta}}\right\rfloor \text {, } \tag{2.1}
\end{equation*}
$$

for a constant $k_{A}>0$. The prior density for $\theta \mid k$ satisfies $\operatorname{Hyp}\left(\left\{\mathrm{k}_{\mathrm{n}}\right\}, \mathrm{c}_{0}, \beta-\right.$ $\frac{1}{2}, \mathrm{~L}_{\mathrm{o}}$ ) for some $c_{0}>0$ and has support $\Theta_{k}\left(\beta-\frac{1}{2}, L\right)$. In addition, for all $\theta, \theta^{\prime} \in \Theta_{k}\left(\beta-\frac{1}{2}, L\right)$ such that $\left\|\theta-\theta^{\prime}\right\| \leq L(n / \log n)^{-\frac{2 \beta-1}{4 \beta}}$,

$$
\begin{equation*}
\log \pi_{\theta \mid k}(\theta)-\log \pi_{\theta \mid k}\left(\theta^{\prime}\right)=h_{k}^{t}\left(\theta-\theta^{\prime}\right)+o(1) \tag{2.2}
\end{equation*}
$$

for constants $C, \rho_{0}>0$ and vectors $h_{k}$ satisfying $\left\|h_{k}\right\| \leq C(n / k)^{1-\rho_{0}}$. Finally, it is assumed that $L$ is sufficiently large compared to $L_{o}$.

- Prior B: $k$ is deterministic and increasing at rate

$$
k_{n}^{\prime}=\left\lfloor k_{B}(n / \log n)^{\frac{1}{1+2 \beta}}\right\rfloor,
$$

where $k_{B}$ is such that $k_{n}^{\prime}<k_{n}$ for all $n$. The prior for $\theta \mid k$ has density $\pi_{\theta \mid k}$ with respect to Lebesgue measure which satisfies condition $\operatorname{Hyp}\left(\left\{\mathrm{k}_{\mathrm{n}}^{\prime}\right\}, \mathrm{c}_{0}, \beta, \mathrm{~L}_{\mathrm{o}}\right)$ for some $c_{0}>0$ and is assumed to have support $\Theta_{k}(\beta, L)$. The density also satisfies

$$
\log \pi_{\theta \mid k}(\theta)-\log \pi_{\theta \mid k}\left(\theta^{\prime}\right)=o(1)
$$

for all $\theta, \theta^{\prime} \in \Theta_{k}(\beta, L)$ such that $\left\|\theta-\theta^{\prime}\right\| \leq L(n / \log n)^{-\frac{\beta}{2 \beta+1}}$. This condition is similar to (2.2), but with $h_{k}=0$, and support $\Theta_{k}(\beta, L)$.

- Prior C: $k \sim \pi_{k}$ on $\mathbb{N}$ with $e^{-c_{1} k \log k} \leq \pi_{k}(k) \leq e^{-c_{2} k \log k}$ for $k$ large enough, where $0<c_{1}<c_{2}<+\infty$. There exists $\beta_{s}>1$ such that for all $\beta \geq \beta_{s}$, the prior for $\theta \mid k$ has density $\pi_{\theta \mid k}$ with respect to Lebesgue measure which satisfies condition $\operatorname{Hyp}\left(\left\{\mathrm{k} \leq \mathrm{k}_{0}(\mathrm{n} / \log \mathrm{n})^{1 /(2 \beta+1)}\right\}, \mathrm{c}_{0}, \beta, \mathrm{~L}_{\mathrm{o}}\right)$, for all $k_{0}>0$ and some $c_{0}>0$, as soon as $n$ is large enough. It has support included in $\Theta_{k}(\beta, L)$ and satisfies

$$
\log \pi_{\theta \mid k}(\theta)-\log \pi_{\theta \mid k}\left(\theta^{\prime}\right)=o(1)
$$

for all $\theta, \theta^{\prime} \in \Theta_{k}(\beta, L)$ such that $\left\|\theta-\theta^{\prime}\right\| \leq L(n / \log n)^{-\frac{\beta}{2 \beta+1}}$.
Note that prior A is obtained when we take $\beta^{\prime}=\beta-\frac{1}{2}$ in prior B.
2.2. Examples of priors. The Lipschitz conditions on $\log \pi_{\theta \mid k}$ considered for the three types of priors are satisfied for instance for the uniform prior on
$\Theta_{k}\left(\beta-\frac{1}{2}, L\right)$ (resp. $\left.\Theta_{k}(\beta, L)\right)$, and for the truncated Gaussian prior, where, for some constants $A$ and $\alpha>0$,

$$
\pi_{\theta \mid k}(\theta) \propto \mathbb{I}_{\Theta_{k}\left(\beta-\frac{1}{2}, L\right)}(\theta) \exp \left(-A \sum_{j=0}^{k} j^{\alpha} \theta_{j}^{2}\right) .
$$

In the case of Prior A, the conditions on $\log \pi_{\theta \mid k}$ and $h_{k}$ in (2.2) are satisfied for $\alpha<4 \beta-2$. To see this, note that for all $\theta, \theta^{\prime} \in \Theta_{h}(\beta-1 / 2, L)$,

$$
\sum_{j=0}^{k} j^{\alpha}\left|\theta_{j}^{2}-\left(\theta_{j}^{\prime}\right)^{2}\right| \leq L^{1 / 2}\left\|\theta-\theta^{\prime}\right\| k^{\alpha-\beta+1 / 2}=o\left((n / k)^{1-\delta}\right)
$$

In the case of Prior B and and Prior C we may choose $\alpha<2 \beta$, since for some positive $k_{0}$

$$
\sum_{j=0}^{k} j^{\alpha}\left|\theta_{j}^{2}-\left(\theta_{j}^{\prime}\right)^{2}\right| \leq L^{1 / 2}\left\|\theta-\theta^{\prime}\right\| k^{\alpha-\beta}=o(1)
$$

for all $k \leq k_{0}(n / \log n)^{1 /(2 \beta+1)}$ and all $\theta, \theta^{\prime} \in \Theta_{k}(\beta, L)$ such that $\left\|\theta-\theta^{\prime}\right\| \leq$ $(n / \log n)^{-\beta /(2 \beta+1)}$.

Also a truncated Laplace distribution is possible, in which case

$$
\pi_{\theta \mid k}(\theta) \propto \mathbb{I}_{\Theta_{k}\left(\beta-\frac{1}{2}, L\right)}(\theta) \exp \left(-a \sum_{j=0}^{k}\left|\theta_{j}\right|\right) .
$$

The condition on $\pi_{k}$ in Prior C is satisfied for instance by Poisson distributions.

The restriction of the prior to Sobolev balls is required to obtain a proper concentration rate or even consistency of the posterior of the spectral density $f$ itself, which is a necessary step in the proof of our results. This is discussed in more detail in section 3.1.
2.3. Convergence rates and BVM-results under different priors. Assuming a Poisson prior for $k$, RCL (Theorem 4.2) obtain a near-optimal convergence rate for $l\left(f, f_{o}\right)$. In Corollary 3.1 below, we show that the optimal rate for $l$ implies that we have at least a suboptimal rate for $\left|d-d_{o}\right|$. Whether this can be improved to the optimal rate critically depends on the prior on $k$. By our first main result the answer is positive under prior A. The proof is given in section 4.

Theorem 2.1. Under prior $\boldsymbol{A}$, the posterior distribution has the asymptotic expansion

$$
\begin{equation*}
\Pi\left[\left.\sqrt{\frac{n r_{k_{n}}}{2}}\left(d-d_{o}-b_{n}\left(d_{o}\right)\right) \leq z \right\rvert\, X\right]=\Phi(z)+o_{P_{o}}(1) \tag{2.3}
\end{equation*}
$$

where, for $r_{k_{n}}=\sum_{j \geq k_{n}+1} \eta_{j}^{2}$ and some small enough $\delta>0$,

$$
b_{n}\left(d_{o}\right)=\frac{1}{r_{k_{n}}} \sum_{j=k_{n}+1}^{\infty} \eta_{j} \theta_{o, j}+Y_{n}+o\left(n^{-1 / 2-\delta} k_{n}^{1 / 2}\right), \quad Y_{n}=\frac{\sqrt{2}}{\sqrt{n r_{k_{n}}}} Z_{n},
$$

$Z_{n}$ being a sequence of random variables converging weakly to a Gaussian variable with mean zero and variance 1 .

Corollary 2.1. Under prior A, the convergence rate for $d$ is $\delta_{n}=$ $(n / \log n)^{-\frac{2 \beta-1}{4 \beta}}$, i.e.

$$
\lim _{n \rightarrow \infty} E_{0}^{n}\left[\Pi\left(d:\left|d-d_{o}\right|>\delta_{n} \mid X\right)\right]=0 .
$$

Equation (2.3) is a Bernstein-von Mises type of result: the posterior distribution is asymptotically normal, centered at a point $d_{o}+b_{n}\left(d_{o}\right)$, whose distribution is normal with mean $d_{o}$ and variance $2 /\left(n r_{k_{n}}\right)$. The expressions for the posterior mean and variance give more insight in how the prior for $k$ affects the posterior rate for $d$. The standard deviation of the limiting normal distribution (2.3) is $\sqrt{2 /\left(n r_{k_{n}}\right)}=O\left(n^{-\frac{2 \beta-1}{4 \beta}}(\log n)^{\frac{1}{4 \beta}}\right)$ and $b_{n}\left(d_{o}\right)$ equals

$$
\frac{1}{r_{k_{n}}} \sum_{j=k_{n}+1}^{\infty} \eta_{j} \theta_{o, j}+O_{P_{o}}\left(k_{n}^{\frac{1}{2}} n^{\frac{1}{2}}\right)+o\left(n^{-1 / 2-\delta_{1}} k_{n}^{1 / 2}\right)
$$

From the definition of $\eta_{j}, k_{n}$ and $r_{k_{n}}$ and the assumption on $\theta_{o}$, it follows that

$$
\begin{equation*}
\frac{1}{r_{k_{n}}}\left|\sum_{j=k_{n}+1}^{\infty} \eta_{j} \theta_{o, j}\right| \leq \frac{1}{r_{k_{n}}} \sqrt{\sum_{l>k_{n}} \theta_{o, l}^{2} j^{2 \beta}} \sqrt{\sum_{l>k_{n}} j^{-2 \beta-2}}=o\left(k_{n}^{-\beta+\frac{1}{2}}\right) . \tag{2.4}
\end{equation*}
$$

See also (1.9) in the supplement. Hence, when the constant $k_{A}$ in (2.1) is small enough,

$$
\begin{equation*}
\left|b_{n}\left(d_{o}\right)\right| \leq \delta_{n}, \tag{2.5}
\end{equation*}
$$

and we obtain the $\delta_{n}$-rate of Corollary 2.1. For smaller $k$, the standard deviation is smaller but the bias $b_{n}\left(d_{o}\right)$ is larger. In Theorem 2.2 below it is shown that this indeed leads to a suboptimal rate.

An important consequence of the BVM-result is that posterior credible regions for $d$ (HPD or equal-tails for instance) will also be asymptotic frequentist confidence regions. Consider for instance one-sided credible intervals for $d$ defined by $P^{\pi}\left(d \leq z_{n}(\alpha) \mid X\right)=\alpha$, so that $z_{n}(\alpha)$ is the $\alpha$-th quantile of the posterior distribution of $d$. Equation (2.3) in Theorem 2.1 then implies that

$$
z_{n}(\alpha)=d_{o}+b_{n}\left(d_{o}\right)+\sqrt{\frac{2 k_{n}}{n}} \Phi^{-1}(\alpha)\left(1+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)\right)
$$

As soon as $\sum_{j \geq k_{n}} j^{2 \beta} \theta_{o, j}^{2}=o\left((\log n)^{-1}\right)$, we have that

$$
z_{n}(\alpha)=d_{o}+\sqrt{2 /\left(n r_{k_{n}}\right)} Z_{n}+\sqrt{2 /\left(n r_{k_{n}}\right)} \Phi^{-1}(\alpha)\left(1+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)\right)
$$

and

$$
P_{o}^{n}\left(d_{o} \leq z_{n}(\alpha)\right)=P\left(Z_{n} \leq \Phi^{-1}(\alpha)(1+o(1))\right)=\alpha+o(1)
$$

Similar computations can be made on equal - tail credible intervals or HPD regions for $d$.

Note that in this paper we assume that the smoothness $\beta$ of $f_{o}$ is greater than 1 instead of $1 / 2$, as is required in Moulines and Soulier (2003). This condition is used throughout the proof. Actually had we only assumed that $\beta>3 / 2$, the proof of Theorem 2.1 would have been greatly simplified as many technicalities in the paper come from controlling terms when $1<$ $\beta \leq 3 / 2$. We do not believe that it is possible to weaken this constraint to $\beta>1 / 2$ in our setup.

Our second main result states that if $k$ is increasing at a slower rate than $k_{n}$, the posterior on $d$ concentrates at a suboptimal rate. The proof is given in section 5 .

Theorem 2.2. Given $\beta>5 / 2$, there exists $\theta_{o} \in \Theta\left(\beta, L_{o}\right)$ and a constant $k_{v}>0$ such that under prior $B$ and $C$ defined above,

$$
\Pi\left(\left|d-d_{o}\right|>k_{v} w_{n}(\log n)^{-1} \mid X\right) \xrightarrow{P_{\rho}} 1
$$

with $w_{n}=C_{w}(n / \log n)^{-\frac{2 \beta-1}{4 \beta+2}}$ and $C_{w}=C_{1}\left(L+L_{o}\right)^{\frac{1}{4 \beta}} l_{0}^{\frac{2 \beta-1}{2 \beta}}$.
The constant $C_{w}$ comes from the suboptimal rate for $\left|d-d_{o}\right|$ derived in Corollary 3.1. Theorem 2.2 is proved by considering the vector $\theta_{o}$ defined by $\theta_{o, j}=c_{0} j^{-\left(\beta+\frac{1}{2}\right)}(\log j)^{-1}$, for $j \geq 2$. This vector is close to the boundary of
the Sobolev-ball $\Theta\left(\beta, L_{o}\right)$, in the sense that for all $\beta^{\prime}>\beta, \sum_{j} j^{2 \beta^{\prime}} \theta_{o, j}^{2}=+\infty$. The proof consists in showing that conditionally on $k$, the posterior distribution is asymptotically normal as in (2.3), with $k$ replacing $k_{n}$, and that the posterior distribution concentrates on values of $k$ smaller than $O\left(n^{1 /(2 \beta+1)}\right)$, so that the bias $b_{n}\left(d_{o}\right)$ becomes of order $w_{n}(\log n)^{-1}$. The constraint $\beta>5 / 2$ is used to simplify the computations and is not sharp.

It is interesting to note that similar to the frequentist approach, a key issue is a bias-variance trade-off, which is optimized when $k \sim n^{1 /(2 \beta)}$. This choice of $k$ depends on the smoothness parameter $\beta$, and since it is not of the same order as the optimal values of $k$ for the loss $l\left(f, f^{\prime}\right)$ on the spectral densities, the adaptive (near) minimax Bayesian nonparametric procedure proposed in Rousseau and Kruijer (2011) does not lead to optimal posterior concentration rate for $d$. While it is quite natural to obtain an adaptive (nearly) minimax Bayesian procedure under the loss $l(.,$.$) by choosing a$ random $k$, obtaining an adaptive minimax procedure for $d$ remains an open problem. This dichotomy is found in other semi-parametric Bayesian problems, see for instance Arbel (2010) in the case of the white noise model or Rivoirard and Rousseau (2010) for BVM properties.
3. Decomposing the posterior for $\boldsymbol{d}$. To prove Theorems 2.1 and 2.2 we need to take a closer look at (1.3), to understand how the integration over $\Theta_{k}$ affects the posterior for $d$. We develop $\theta \rightarrow l_{n}(d, k, \theta)$ in a point $\bar{\theta}_{d, k}$ defined below and decompose the likelihood as

$$
\exp \left\{l_{n}(d, k, \theta)\right\}=\exp \left\{l_{n}(d, k)\right\} \exp \left\{l_{n}(d, k, \theta)-l_{n}(d, k)\right\},
$$

where $l_{n}(d, k)$ is short-hand notation for $l_{n}\left(d, k, \bar{\theta}_{d, k}\right)$. Define

$$
\begin{equation*}
I_{n}(d, k)=\int_{\Theta_{k}} e^{l_{n}(d, k, \theta)-l_{n}(d, k)} d \pi_{\theta \mid k}(\theta), \tag{3.1}
\end{equation*}
$$

where $\Theta_{k}$ is the generic notation for $\Theta_{k}\left(\beta-\frac{1}{2}, L\right)$ under prior $\mathbf{A}$ and $\Theta_{k}(\beta, L)$ for priors B and C. The posterior for $d$ given in (1.3) can be written as

$$
\begin{equation*}
\Pi(d \in D \mid X)=\frac{\sum_{k=0}^{\infty} \pi_{k}(k) \int_{D} e^{l_{n}(d, k)-l_{n}\left(d_{o}, k\right)} I_{n}(d, k) d \pi_{d}(d)}{\sum_{k=0}^{\infty} \pi_{k}(k) \int_{-\frac{1}{2}+t}^{\frac{1}{2}-t} e^{l_{n}(d, k)-l_{n}\left(d_{o}, k\right)} I_{n}(d, k) d \pi_{d}(d)} \tag{3.2}
\end{equation*}
$$

The factor $\exp \left\{l_{n}(d, k)-l_{n}\left(d_{o}, k\right)\right\}$ is independent of $\theta$, and will under certain conditions dominate the marginal likelihood. In section 3.2 we give a Taylor-approximation which, for given $k$, allows for a normal approximation to the marginal posterior. However, to obtain the convergence rates in

Theorems 2.1 and 2.2, it also needs to be shown that the integrals $I_{n}(d, k)$ with respect to $\theta$ do not vary too much with $d$. This is the most difficult part of the proof of Theorem 2.1 and the argument is presented in section 3.3. Since Theorem 2.2 is essentially a counter-example and it is not aimed to be as general as Theorem 2.1, as far as the range of $\beta$ is concerned, we can restrict attention to larger $\beta^{\prime}$ 's, i.e. $\beta>5 / 2$, for which controlling $I_{n}(d, k)$ is much easier.
3.1. Preliminaries. First we define the point $\bar{\theta}_{d, k}$ in which we develop $\theta \rightarrow l_{n}(d, k, \theta)$. Since the function $\log (2-2 \cos (x))$ has Fourier coefficients against $\cos j x, j \in \mathbb{N}$ equal to $0,2, \frac{2}{2}, \frac{2}{3}, \ldots$, FEXP-spectral densities can be written as

$$
\left|1-e^{i x}\right|^{-2 d} \exp \left\{\sum_{j=0}^{\infty} \theta_{j} \cos (j x)\right\}=\exp \left\{\sum_{j=0}^{\infty}\left(\theta_{j}+d \eta_{j}\right) \cos (j x)\right\}
$$

Given $f=f_{d, k, \theta}$ and $f^{\prime}=f_{d^{\prime}, k^{\prime}, \theta^{\prime}}$ we can therefore express the norm $l\left(f, f^{\prime}\right)$ in terms of $\left(\theta-\theta^{\prime}\right)$ and $\left(d-d^{\prime}\right)$ :

$$
\begin{equation*}
l\left(f, f^{\prime}\right)=\frac{1}{2} \sum_{j=0}^{\infty}\left(\left(\theta_{j}-\theta_{j}^{\prime}\right)+\eta_{j}\left(d-d^{\prime}\right)\right)^{2} \tag{3.3}
\end{equation*}
$$

where $\theta_{j}$ and $\theta_{j}^{\prime}$ are understood to be zero when $j$ is larger than $k$ respectively $k^{\prime}$. Equation (3.3) implies that for given $d$ and $k, l\left(f_{o}, f_{d, k, \theta}\right)$ is minimized by

$$
\bar{\theta}_{d, k}:=\operatorname{argmin}_{\theta \in \mathbb{R}^{k+1}} \sum_{j=0}^{\infty}\left(\theta_{j}-\theta_{o, j}+\left(d-d_{o}\right) \eta_{j}\right)^{2}=\theta_{o[k]}+\left(d_{o}-d\right) \eta_{[k]}
$$

In particular, $\theta=\theta_{o[k]}$ minimizes $l\left(f_{o}, f_{d, k, \theta}\right)$ only when $d=d_{o}$; when $d \neq d_{o}$ we need to add $\left(d_{o}-d\right) \eta_{[k]}$. The following lemma shows that an upper bound on $l\left(f_{o}, f_{d, k, \theta}\right)$ leads to upper bounds on $\left|d-d_{o}\right|$ and $\left\|\theta-\theta_{o}\right\|$.

Lemma 3.1. Suppose that $\theta \in \Theta_{k}(\gamma, L)$ and $\theta_{o} \in \Theta_{k}\left(\beta, L_{o}\right)$, where $\gamma \leq$ $\beta$. Also suppose that for a sequence $\alpha_{n} \rightarrow 0, l\left(f_{o}, f_{d, k, \theta}\right) \leq \alpha_{n}^{2}$ for all $n$. Then there are universal constants $C_{1}, C_{2}>0$ such that for all n,

$$
\left|d-d_{o}\right| \leq C_{1}\left(L+L_{o}\right)^{\frac{1}{4 \gamma}} \alpha_{n}^{\frac{2 \gamma-1}{2 \gamma}}, \quad\left\|\theta-\theta_{o}\right\| \leq C_{2}\left(L+L_{o}\right)^{\frac{1}{4 \gamma}} \alpha_{n}^{\frac{2 \gamma-1}{2 \gamma}}
$$

Proof. For all $(d, k, \theta)$ such that $l\left(f_{d, k, \theta}, f_{o}\right) \leq \alpha_{n}$, we have, using (3.3),

$$
\begin{aligned}
2 \alpha_{n}^{2} & \geq 2 l\left(f_{d, k, \theta}, f_{o}\right)=2\left(\theta_{o, 0}-\theta_{0}\right)^{2}+\sum_{j \geq 1}\left(\left(\theta_{o, j}-\theta_{j}\right)+\eta_{j}\left(d_{o}-d\right)\right)^{2} \\
& \geq \sum_{j \geq 1}\left(\theta_{o, j}-\theta_{j}\right)^{2}+\left(d-d_{o}\right)^{2} \sum_{j \geq 1} \eta_{j}^{2}-2\left|d-d_{o}\right| \sqrt{\sum_{j \geq 1} \eta_{j}^{2}} \sqrt{\sum_{j \geq 1}\left(\theta_{o, j}-\theta_{j}\right)^{2}} \\
& =\left(\left\|\theta-\theta_{o}\right\|-\left|d-d_{o}\right|\|\eta\|\right)^{2}
\end{aligned}
$$

The inequalities remain true if we replace all sums over $j \geq 1$ by sums over $j \geq m_{n}$, for any nondecreasing sequence $m_{n}$. Since $\left\|\left(\eta_{j} 1_{j>m_{n}}\right)_{j \geq 1}\right\|^{2}$ is of order $m_{n}^{-1}$ and $\left\|\left(\theta-\theta_{o} \|_{j} 1_{j>m_{n}}\right)_{j \geq 1}\right\|^{2} \leq m_{n}^{-2 \gamma} \sum_{j>m_{n}}(1+j)^{2 \beta}\left(\theta_{j}-\theta_{o, j}\right)^{2}<$ $2\left(L+L_{o}\right) m_{n}^{-2 \gamma}$, setting $m_{n}=\alpha_{n}^{-\frac{1}{\gamma}}$ gives the desired rate for $\left|d-d_{o}\right|$ as well as for $\left\|\theta-\theta_{o}\right\|$.

The convergence rate for $l\left(f_{o}, f_{d, k, \theta}\right)$ required in Lemma 3.1 can be found in Rousseau and Kruijer (2011). For easy reference we restate it here. Compared to a similar result in RCL, the $\log n$ factor is improved.

Lemma 3.2. Under prior $\boldsymbol{A}$, there exists a constant $l_{0}$ depending only on $L_{o}$ and $k_{A}$ (and not on $L$ ) such that

$$
\Pi\left((d, k, \theta): l\left(f_{d, k, \theta}, f_{o}\right) \geq l_{0}^{2} \delta_{n}^{2} \mid X\right) \xrightarrow{P_{\rho}} 0
$$

where $\delta_{n}=(n / \log n)^{-\frac{2 \beta-1}{4 \beta}}$. Under priors $\boldsymbol{B}$ and $\mathbf{C}$, this statement holds with $\epsilon_{n}=(n / \log n)^{-\frac{\beta}{2 \beta+1}}$ replacing $\delta_{n}$.

In the proof of Theorem 2.1 (resp. 2.2), this result allows us to restrict attention to the set of spectral densities $f$ such that $l\left(f, f_{o}\right) \leq l_{0}^{2} \delta_{n}^{2}$ (resp. $\left.l_{0}^{2} \epsilon_{n}^{2}\right)$. In addition, by combination with Lemma 3.1 we can now deduce bounds on $\left|d-d_{o}\right|$ and $\left\|\theta-\bar{\theta}_{d, k}\right\|$. These bounds, although suboptimal, will be important in the sequel for obtaining the near-optimal rate in Theorem 2.1.

Corollary 3.1. Under the result of Lemma 3.2 and prior A, we can apply Lemma 3.1 with $\alpha_{n}^{2}=l_{0}^{2} \delta_{n}^{2}$ and $\gamma=\beta-\frac{1}{2}$, and obtain

$$
\Pi_{d}\left(d:\left|d-d_{o}\right| \geq \bar{v}_{n} \mid X\right) \xrightarrow{P_{o}} 0, \quad \Pi\left(\left\|\theta-\bar{\theta}_{d, k}\right\| \geq 2 l_{0} \delta_{n} \mid X\right) \xrightarrow{P_{\rho}} 0
$$

where $\bar{v}_{n}=C_{1}\left(L+L_{o}\right)^{\frac{1}{4 \beta-2}} l_{0}^{\frac{2 \beta-2}{2 \beta-1}}(n / \log n)^{-\frac{\beta-1}{2 \beta}}$. Under priors $\boldsymbol{B}$ and $\mathbf{C}$ we have $\gamma=\beta$; the rate for $\left|d-d_{o}\right|$ is then $w_{n}=C_{w}(n / \log n)^{-\frac{2 \beta-1}{4 \beta+2}}$ and the
rate for $\left\|\theta-\bar{\theta}_{d, k}\right\|$ is $2 l_{0} \epsilon_{n}$. The constant $C_{w}=C_{1}\left(L+L_{o}\right)^{\frac{1}{4 \beta}} l_{0}^{\frac{2 \beta-1}{2 \beta}}$ is as in Theorem 2.2.

Proof. The rate for $\left|d-d_{o}\right|$ follows directly from Lemma 3.1. To obtain the rate for $\left\|\theta-\bar{\theta}_{d, k}\right\|$, let $\alpha_{n}$ denote either $l_{0} \delta_{n}$ (the rate for $l\left(f_{o}, f\right)$ under prior A) or $l_{0} \epsilon_{n}$ (the rate under priors B and C ). Although Lemma 3.1 suggests that the Euclidean distance from $\theta_{o}$ to $\theta$ (contained in $\Theta_{k}(\beta, L)$ or $\left.\Theta_{k}\left(\beta-\frac{1}{2}, L\right)\right)$ may be larger than $\alpha_{n}$, the distance from $\theta$ to $\bar{\theta}_{d, k}$ is certainly of order $\alpha_{n}$. To see this, note that Lemma 3.2 implies the existence of $d, k, \theta$ in the model with $l\left(f_{o}, f_{d, k, \theta}\right) \leq \alpha_{n}^{2}$. From the definition of $\bar{\theta}_{d, k}$ it follows that $l\left(f_{o}, f_{d, k, \bar{\theta}_{d, k}}\right) \leq \alpha_{n}^{2}$. The triangle inequality gives $\left\|\theta-\bar{\theta}_{d, k}\right\|^{2}=$ $l\left(f_{d, k, \theta}, f_{d, k, \bar{\theta}_{d, k}}\right) \leq 4 \alpha_{n}^{2}$.

The rates $\bar{v}_{n}$ and $w_{n}$ obtained in Corollary 3.1 are clearly suboptimal; their importance however lies in the fact that they narrow down the set for which we need to prove Theorems 2.1 and 2.2. To prove Theorem 2.2 for example it suffices to show that the posterior mass on $k_{v} w_{n}(\log n)^{-1}<\left|d-d_{o}\right|<w_{n}$ tends to zero. Note that the lower and the upper bound differ only by a factor $(\log n)$. Hence under priors B and C, the combination of Corollary 3.1 and Theorem 2.2 characterizes the posterior concentration rate (up to a $\log n$ term) for the given $\theta_{o}$. Another consequence of Corollary 3.1 is that we may neglect the posterior mass on all $(d, k, \theta)$ for which $\left\|\theta-\bar{\theta}_{d, k}\right\|$ is larger than $2 l_{0} \delta_{n}$ (under prior A) or $2 l_{0} \epsilon_{n}$ (under priors B and C).

We conclude this section with a result on $\bar{\theta}_{d, k}$ and $\Theta_{k}(\beta, L)$. In the definition of $\bar{\theta}_{d, k}$ we minimize over $\mathbb{R}^{k+1}$, whereas the support of priors A-C is the Sobolev ball $\Theta_{k}(\beta, L)$ or $\Theta_{k}\left(\beta-\frac{1}{2}, L\right)$. Under the assumptions of Theorems 2.1 and 2.2 however, $\bar{\theta}_{d, k}$ is contained in $\Theta_{k}\left(\beta-\frac{1}{2}, L\right)$ respectively $\Theta_{k}(\beta, L)$. Also the $l_{2}$-ball of radius $2 l_{0} \delta_{n}$ (or $2 l_{0} \epsilon_{n}$ ) is contained in these Sobolev-balls.

Lemma 3.3. Under the assumptions of Theorem 2.1, $B_{k}\left(\bar{\theta}_{d, k}, 2 l_{0} \delta_{n}\right)$ is contained in $\Theta_{k}\left(\beta-\frac{1}{2}, L\right)$ for all $d \in\left[d_{o}-\bar{v}_{n}, d_{o}+\bar{v}_{n}\right]$, if $L$ is large enough. In particular, $\bar{\theta}_{d, k} \in \Theta_{k}\left(\beta-\frac{1}{2}, L\right)$. Similarly, under the assumptions of Theorem 2.2, $B_{k}\left(\bar{\theta}_{d, k}, 2 l_{0} \epsilon_{n}\right) \subset \Theta_{k}(\beta, L)$, for all $d \in\left[d_{o}-w_{n}, d_{o}+w_{n}\right]$.

Proof. Since the constant $l_{0}$ is independent of $L, \theta \in B_{k}\left(\bar{\theta}_{d, k}, 2 l_{0} \delta_{n}\right)$
implies that for $n$ large enough ,

$$
\begin{aligned}
& \sum_{j=0}^{k} \theta_{j}^{2}(j+1)^{2 \beta-1} \leq 2 \sum_{j=0}^{k}\left(\theta-\bar{\theta}_{d, k}\right)_{j}^{2}(j+1)^{2 \beta-1}+2 \sum_{j=0}^{k}\left(\bar{\theta}_{d, k}\right)_{j}^{2}(j+1)^{2 \beta-1} \\
& \leq \\
& \quad 8 \delta^{2}\left(L_{o}\right)(n / \log n)^{\frac{2 \beta-1}{2 \beta}}\left(k_{n}+1\right)^{2 \beta-1}+4 \sum_{j=0}^{k_{n}} \theta_{o, j}^{2}(j+1)^{2 \beta-1} \\
& \quad+16\left(d-d_{o}\right)^{2} \sum_{j=1}^{k_{n}} j^{2 \beta-3} .
\end{aligned}
$$

The first two terms on the right only depend on $L_{o}$, and are smaller than $L / 4$ when $L$ is chosen sufficiently large. Because $\bar{v}_{n}=C_{1}\left(L+L_{o}\right)^{\frac{1}{\beta \beta-2}} l_{0}^{\frac{2 \beta-2}{2 \beta-1}}(n / \log n)^{-\frac{\beta-1}{2 \beta}}$, the last term in the preceding display is at most

$$
C_{1}^{2}\left(L+L_{o}\right)^{\frac{1}{2 \beta-1}} l_{0}^{\frac{4 \beta-4}{2-1}}(n / \log n)^{-\frac{\beta-1}{\beta}} k_{A}^{2 \beta-2}(n / \log n)^{\frac{\beta-1}{\beta}},
$$

which, since $\beta>1$, is smaller than $L / 2$ when $L$ is large enough. We conclude that $B_{k}\left(\bar{\theta}_{d, k}, 2 l_{0} \delta_{n}\right)$ is contained in $\Theta_{k}\left(\beta-\frac{1}{2}, L\right)$ provided $L$ is chosen sufficiently large. The second statement can be proved similarly.
3.2. A Taylor approximation for $l_{n}(d, k)$. Provided that the integrals $I_{n}(d, k)$ have negligible impact on the posterior for $d$, the conditional distribution of $d$ given $k$ will only depend on $\exp \left\{l_{n}(d, k)-l_{n}\left(d_{o}, k\right)\right\}$. Let $l_{n}^{(1)}(d, k), l_{n}^{(2)}(d, k)$ denote the first two derivatives of the map $d \mapsto l_{n}(d, k)$. There exists a $\bar{d}$ between $d$ and $d_{o}$ such that

$$
\begin{equation*}
l_{n}(d, k)=l_{n}\left(d_{o}, k\right)+\left(d-d_{o}\right) l_{n}^{(1)}\left(d_{o}, k\right)+\frac{\left(d-d_{o}\right)^{2}}{2} l_{n}^{(2)}(\bar{d}, k) \tag{3.4}
\end{equation*}
$$

Defining

$$
b_{n}(d)=-\frac{l_{n}^{(1)}\left(d_{o}, k\right)}{l_{n}^{(2)}(d, k)}
$$

which is the $b_{n}$ used in Theorem 2.1, we can rewrite (3.4) as

$$
\begin{equation*}
l_{n}(d, k)-l_{n}\left(d_{o}, k\right)=-\frac{1}{2} \frac{\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)^{2}}{l_{n}^{(2)}(\bar{d}, k)}+\frac{1}{2} l_{n}^{(2)}(\bar{d}, k)\left(d-d_{o}-b_{n}(\bar{d})\right)^{2} . \tag{3.5}
\end{equation*}
$$

Note that each derivative $l_{n}^{(i)}(d, k), i=1,2$, can be decomposed into a centered quadratic form denoted $\mathcal{S}\left(l_{n}^{(i)}(d, k)\right)$ and a deterministic term $\mathcal{D}\left(l_{n}^{(i)}(d, k)\right)$.

In the following lemma we give expressions for $l_{n}^{(1)}\left(d_{o}, k\right), l_{n}^{(2)}(d, k)$ and $b_{n}$, making explicit their dependence on $k$ and $\theta_{o}$. Since $k_{n}^{\prime} \leq k_{n}$ and $w_{n}<\bar{v}_{n}$ (see Corollary 3.1) the result is valid for all priors under consideration. The proof is given in appendix A.

Lemma 3.4. Given $\beta>1$, let $\theta_{o} \in \Theta\left(\beta, L_{o}\right)$. If $k \leq k_{n}$ and $\left|d-d_{o}\right| \leq \bar{v}_{n}$, then there exists $\delta_{1}>0$ such that

$$
\begin{aligned}
l_{n}^{(1)}\left(d_{o}, k\right) & :=\mathcal{S}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)+\mathcal{D}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right) \\
& =\mathcal{S}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)+\frac{n}{2} \sum_{j=k+1}^{\infty} \theta_{o, j} \eta_{j}+o\left(n^{\epsilon}\left(k^{-\beta+3 / 2}+n^{-1 /(2 \beta)}\right)\right) \\
l_{n}^{(2)}(d, k) & =l_{n}^{(2)}\left(d_{o}, k\right)\left(1+\frac{k^{1 / 2}}{n^{1 / 2+\epsilon}}+\frac{k^{-2 \beta+1+\epsilon}}{n}\right)=-\frac{1}{2} n r_{k}\left(1+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{-\delta_{1}}\right)\right),
\end{aligned}
$$

where $\mathcal{S}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)$ is a centered quadratic form with variance

$$
\operatorname{Var}\left(\mathcal{S}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)\right)=\frac{n}{2} \sum_{j>k} \eta_{j}^{2}(1+o(1))=\frac{n r_{k}}{2}(1+o(1))=O\left(n k^{-1}\right) .
$$

Consequently,

$$
\begin{align*}
b_{n}(d)=- & \frac{l_{n}^{(1)}\left(d_{o}, k\right)}{l_{n}^{(2)}(d, k)}=\frac{1}{r_{k}} \sum_{j=k+1}^{\infty} \theta_{o, j} \eta_{j}\left(1+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{-\delta}\right)\right)  \tag{3.6}\\
& +\frac{2 \mathcal{S}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)\left(1+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{-\delta}\right)\right)}{n r_{k}}+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{\epsilon-1} k^{-\beta+5 / 2}+n^{\epsilon-1}\right),
\end{align*}
$$

with

$$
\frac{2 \mathcal{S}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)}{n r_{k}}=\mathbf{O}_{\mathbf{P}_{\mathbf{o}}}\left(n^{-\frac{1}{2}} k^{\frac{1}{2}}\right)
$$

Remark 3.1. Recall from (2.4) that $r_{k}^{-1} \sum_{j=k+1}^{\infty} \theta_{o, j} \eta_{j}$ is $O\left(k^{-\beta+1 / 2}\right)$. The term $2 \mathcal{S}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right) /\left(n r_{k}\right)$ is $O_{P_{o}}\left(k^{-\beta+1 / 2}\right)$ whenever $k \sim n^{1 /(2 \beta)}$, which is the case under all priors under consideration.

Substituting the above results on $l_{n}^{(1)}, l_{n}^{(2)}$ and $b_{n}$ in (3.5), we can give the following informal argument leading to Theorems 2.1 and Theorem 2.2. If we consider $k$ to be fixed and $I_{n}(d, k)$ constant in $d$, then (3.5) implies that the posterior distribution for $d$ is asymptotically normal with mean $d_{o}+b_{n}\left(d_{o}\right)$ and variance of order $k / n$.
3.3. Integration of the short memory parameter. A key ingredient in the proofs of both Theorems 2.1 and 2.2 is the control of the integral $I_{n}(d, k)$ appearing in (1.3), whose dependence on $d$ should be negligible with respect to $\exp \left\{l_{n}(d, k)-l_{n}\left(d_{o}, k\right)\right\}$. In Lemma 3.5 below we prove this to be the case under the assumptions of Theorems 2.1 and 2.2. For the case of Theorem 2.2 this is fairly simple: the conditional posterior distribution of $\theta$ given $(d, k)$ can be proved to be asymptotically Gaussian by a Laplace-approximation. For smaller $\beta$ and larger $k$ the control is technically more demanding. In both cases the proof is based on the following Taylor expansion of $l_{n}(d, k, \theta)$ around $\bar{\theta}_{d, k}$ :

$$
\begin{equation*}
l_{n}(d, k, \theta)-l_{n}(d, k)=\sum_{j=1}^{J} \frac{\left(\theta-\bar{\theta}_{d, k}\right)^{(j)} \nabla^{j} l_{n}(d, k)}{j!}+R_{J+1, d}(\theta) \tag{3.7}
\end{equation*}
$$

where

$$
\begin{align*}
& \quad\left(\theta-\bar{\theta}_{d, k}\right)^{(j)} \nabla^{j} l_{n}(d, k)=\sum_{l_{1}, \ldots, l_{j}=0}^{k}\left(\theta-\bar{\theta}_{d, k}\right)_{l_{1}} \ldots\left(\theta-\bar{\theta}_{d, k}\right)_{l_{j}} \frac{\partial^{j} l_{n}\left(d, k, \bar{\theta}_{d, k}\right)}{\partial \theta_{l_{1}} \ldots \partial \theta_{l_{j}}} \\
& (3.8)  \tag{3.8}\\
& R_{J+1, d}(\theta)=\frac{1}{(J+1)!} \sum_{l_{1}, \ldots, l_{J+1}=0}^{k}\left(\theta-\bar{\theta}_{d, k}\right)_{l_{1}} \ldots\left(\theta-\bar{\theta}_{d, k}\right)_{l_{J+1}} \frac{\partial^{J+1} l_{n}(d, k, \tilde{\theta})}{\partial \theta_{l_{1}} \ldots \partial \theta_{l_{J+1}}} .
\end{align*}
$$

The above expressions are used to derive the following lemma, which gives control of the term $I_{n}(d, k)$.

Lemma 3.5. Under the conditions of Theorem 2.1, the integral $I_{n}(d, k)$ defined in (3.1) equals

$$
I_{n}\left(d_{o}, k\right) \exp \left\{\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(\frac{\left|d-d_{o}\right| n^{\frac{1}{2}-\delta_{2}}}{\sqrt{k}}\right)+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(\left(d-d_{o}\right)^{2} \frac{n^{1-\delta_{2}}}{k}\right)\right\}
$$

for some $\delta_{2}>0$. Under the conditions of Theorem 2.2,

$$
I_{n}(d, k)=I_{n}\left(d_{o}, k\right) \exp \left\{\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)\right\}
$$

The proof is given in Appendix C, and relies on the expressions for the derivatives $\nabla^{j} l_{n}$ given in Appendix B. Lemma 3.5 should be seen in relation to Lemma 3.4 and the expressions for $\Pi(d \mid X)$ and $l_{n}(d, k)-l_{n}\left(d_{o}, k\right)$ in equations (3.2) and (3.4). Lemma 3.5 then shows that the dependence on the integrals $I_{n}(d, k)$ on $d$ is asymptotically negligible with respect to $l_{n}(d, k)-$ $l_{n}\left(d_{o}, k\right)$. This is made rigorous in the following section.
4. Proof of Theorem 2.1. By Lemma 3.2 we may assume posterior convergence of $l\left(f_{o}, f_{d, k, \theta}\right)$ at rate $l_{0}^{2} \delta_{n}^{2}$, and, by Corollary 3.1, also convergence of $\left|d-d_{o}\right|$ at rate $\bar{v}_{n}$. By Lemma 3.3, we may restrict the integration over $\theta$ to $B_{k}\left(\bar{\theta}_{d, k}, 2 l_{0} \delta_{n}\right)$. Let $\Gamma_{n}(z)=\left\{d: \sqrt{\frac{n r_{k}}{2}}\left(d-d_{o}-b_{n}\left(d_{o}\right)\right) \leq z\right\}$. Under prior A, it suffices to show that for $k=k_{n}$,

$$
\begin{align*}
\frac{N_{n}}{D_{n}} & :=\frac{\int_{\Gamma_{n}(z)} e^{l_{n}(d, k)-l_{n}\left(d_{o}, k\right)} \int_{B_{k}\left(\bar{\theta}_{d, k}, 2 l_{0} \delta_{n}\right)} e^{l_{n}(d, k, \theta)-l_{n}(d, k)} d \pi_{\theta \mid k}(\theta) d \pi_{d}(d)}{\int_{\left|d-d_{o}\right|<\bar{v}_{n}} e^{l_{n}(d, k)-l_{n}\left(d_{o}, k\right)} \int_{B_{k}\left(\bar{\theta}_{d, k}, 2 l_{0} \delta_{n}\right)} e^{l_{n}(d, k, \theta)-l_{n}(d, k)} d \pi_{\theta \mid k}(\theta) d \pi_{d}(d)}  \tag{4.1}\\
& =\frac{\int_{\Gamma_{n}(z)} \exp \left\{l_{n}(d, k)-l_{n}\left(d_{o}, k\right)+\log I_{n}(d, k)\right\} d \pi_{d}(d)}{\int_{\left|d-d_{o}\right|<\bar{v}_{n}} \exp \left\{l_{n}(d, k)-l_{n}\left(d_{o}, k\right)+\log I_{n}(d, k)\right\} d \pi_{d}(d)}=\Phi(z)+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)
\end{align*}
$$

Using the results for $l_{n}(d, k)-l_{n}\left(d_{o}, k\right)$ and $I_{n}(d, k)$ given by Lemmas 3.4 and 3.5, we show that for $A_{n} \subset \mathbb{R}^{n}$ defined below such that $P_{o}^{n}\left(A_{n}\right) \rightarrow 1$,

$$
\begin{equation*}
\frac{N_{n}}{D_{n}} \leq \Phi(z)+o(1), \quad \frac{N_{n}}{D_{n}} \geq \Phi(z)+o(1), \quad \forall X \in A_{n} \tag{4.2}
\end{equation*}
$$

Since $P_{o}^{n}\left(A_{n}\right) \rightarrow 1$ this implies the last equality in (4.1).
Note that Lemmas 3.4 and 3.5 also hold for all $\delta_{1}^{\prime}<\delta_{1}$ and $\delta_{2}^{\prime}<\delta_{2}$. In the remainder of the proof, let $0<\delta \leq \min \left(\delta_{1}, \delta_{2}\right)$. For notational simplicity, let $\mathcal{D}=\mathcal{D}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right.$, the deterministic part of $l_{n}^{(1)}\left(d_{o}, k\right)$. For a sufficiently large constant $C_{1}$ and arbitrary $\epsilon_{1}>0$, let $A_{n}$ be the set of $X \in \mathbb{R}^{n}$ such that

$$
\left.\begin{array}{l}
\left|\log I_{n}(d, k)-\log I_{n}\left(d_{o}, k\right)\right| \leq \epsilon_{1}+\left(d-d_{o}\right)^{2} k^{-1} n^{1-\delta}+\left|d-d_{o}\right| k^{-\frac{1}{2}} n^{\frac{1}{2}-\delta} \\
\left|l_{n}^{(1)}\left(d_{o}, k\right)-\mathcal{D}\right| \leq C_{1} n^{\frac{1}{2}} k^{-\frac{1}{2}} \sqrt{\log n}, \quad\left|l_{n}^{(2)}(d, k)+\frac{1}{2} n r_{k}\right| \leq n^{1-\delta} k^{-1}
\end{array}\right\}
$$

for all $\left|d-d_{o}\right| \leq \bar{v}_{n}$. Since $k=k_{n}$ and $\beta>1$, Lemmas 3.4 and 3.5 imply that $P_{o}^{n}\left(A_{n}^{c}\right) \rightarrow 0$. We prove the first inequality in (4.2); the second one can be obtained in the same way. Using (3.4) and the definition of $A_{n}$, it follows
that for all $X \in A_{n}$,

$$
\begin{align*}
& l_{n}(d, k)-l_{n}\left(d_{o}, k\right)+\log I_{n}(d, k)-\log I_{n}\left(d_{o}, k\right) \leq \epsilon_{1}+\left(d-d_{o}\right)^{2} \frac{n^{1-\delta}}{k}  \tag{4.3}\\
& \quad+\left|d-d_{o}\right| \frac{n^{\frac{1}{2}-\delta}}{k^{\frac{1}{2}}}+\left(d-d_{o}\right) l_{n}^{(1)}\left(d_{o}, k\right)-\frac{n r_{k}}{4}\left(d-d_{o}\right)^{2}\left(1-n^{-\delta}\right) \\
& \leq 2 \epsilon_{1}-\frac{n r_{k}}{4}\left(1-\frac{2}{n^{\delta}}\right)\left(d-d_{o}-\frac{2 l_{n}^{(1)}\left(d_{o}, k\right)}{\left(1-\frac{2}{n^{\delta}}\right) n r_{k}}\right)^{2}+\left|d-d_{o}\right| \frac{n^{\frac{1}{2}-\delta}}{k^{\frac{1}{2}}}+\frac{\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)^{2}}{\left(1-\frac{2}{n^{\delta}}\right) n r_{k}} \\
& \leq 3 \epsilon_{1}-\frac{n r_{k}}{4}\left(1-\frac{2}{n^{\delta}}\right)\left(d-d_{o}-\frac{b_{n}\left(d_{o}, k\right)}{1-\frac{2}{n^{\delta}}}\right)^{2} \\
& \quad+\left|d-d_{o}-\frac{b_{n}\left(d_{o}, k\right)}{1-\frac{2}{n^{\delta}}}\right| \frac{n^{\frac{1}{2}-\delta}}{k^{\frac{1}{2}}}+\frac{\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)^{2}}{\left(1-\frac{2}{n^{\delta}}\right) n r_{k}}
\end{align*}
$$

The third inequality follows from (2.5) and Remark 3.1, by which $b_{n}\left(d_{o}\right)=$ $O\left(k^{-\beta+\frac{1}{2}}\right)=O\left(\delta_{n}\right)$. This implies that $\left|b_{n}\left(d_{o}\right)\right| k^{-\frac{1}{2}} n^{\frac{1}{2}-\delta}<\epsilon_{1}$, again for large enough $n$. Similar to the preceding display, we have the lower-bound

$$
\begin{align*}
& l_{n}(d, k)-l_{n}\left(d_{o}, k\right)+\log I_{n}(d, k)-\log I_{n}\left(d_{o}, k\right) \\
& \quad \geq-3 \epsilon_{1}-\frac{n r_{k}}{4}\left(1+2 n^{-\delta}\right)\left(d-d_{o}-\frac{b_{n}\left(d_{o}, k\right)}{\left(1+2 n^{-\delta}\right)}\right)^{2}  \tag{4.4}\\
& \quad-\left|d-d_{o}-\frac{b_{n}\left(d_{o}, k\right)}{\left(1+2 n^{-\delta}\right)}\right| k^{-\frac{1}{2}} n^{\frac{1}{2}-\delta}+\frac{\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)^{2}}{\left(1+2 n^{-\delta}\right) n r_{k}} .
\end{align*}
$$

Note that

$$
\begin{equation*}
\exp \left\{\frac{\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)^{2}}{\left(1-2 n^{-\delta}\right) n r_{k}}-\frac{\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)^{2}}{\left(1+2 n^{-\delta}\right) n r_{k}}\right\}=\exp \{o(1)\} \tag{4.5}
\end{equation*}
$$

which follows from the expression for $l_{n}^{(1)}\left(d_{o}, k\right)$ in Lemma 3.4, the definition of $A_{n}$ and the assumption that $X \in A_{n}$. Therefore, substituting (4.3) in $N_{n}$ and (4.4) in $D_{n}$, the terms $\frac{\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)^{2}}{4 n r_{k}}$ cancel out and by (4.5) we can neglect the difference between $\frac{\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)^{2}}{\left(1 \pm 2 n^{-\delta}\right) n r_{k}}$ and $\frac{\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)^{2}}{n r_{k}}$.

To conclude the proof that $N_{n} / D_{n} \leq \Phi(z)+o(1)$ for each $X \in A_{n}$, we make the change of variables

$$
u=\sqrt{\frac{n r_{k}}{2}\left(1 \pm 2 n^{-\delta}\right)}\left(d-d_{o}-\frac{b_{n}\left(d_{o}\right)}{1 \pm 2 n^{-\delta}}\right),
$$

where we take + in the lower bound for $D_{n}$ and - in the upper-bound for $N_{n}$. Using once more that $b_{n}\left(d_{o}\right)=O\left(\delta_{n}\right)$, we find that for large enough $n,|u| \leq \frac{\bar{v}_{n}}{4} \sqrt{n r_{k}}$ implies $\left|d-d_{o}\right| \leq \bar{v}_{n}$. Hence we may integrate over $|u| \leq$ $\frac{\bar{v}_{n}}{4} \sqrt{n r_{k}}$ in the lower-bound for $D_{n}$. In the upper-bound for $N_{n}$ we may integrate over $u \leq z+\epsilon_{1}$.

Combining (4.3)-(4.5), it follows that for all $\epsilon_{1}$ and all $X \in A_{n}$,

$$
\begin{aligned}
& \frac{N_{n}}{D_{n}} \leq e^{7 \epsilon_{1}}\left(\frac{1+2 n^{-\delta}}{1-2 n^{-\delta}}\right)^{\frac{1}{2}} \frac{\int_{u<z+\epsilon_{1}} \exp \left\{-\frac{1}{2} u^{2}+C n^{-\delta}|u|\right\} d u}{\int_{|u| \leq \leq \frac{\bar{v}_{n}}{4} \sqrt{n r_{k}}} \exp \left\{-\frac{1}{2} u^{2}-C n^{-\delta}|u|\right\} d u} \\
& \leq e^{8 \epsilon_{1}} \frac{\int_{u<z+\epsilon_{1}} \exp \left\{-\frac{1}{2} u^{2}+C n^{-\delta}|u|\right\} d u}{\int_{|u| \leq \frac{\bar{v}_{n}}{4} \sqrt{n r_{k}}} \exp \left\{-\frac{1}{2} u^{2}-C n^{-\delta}|u|\right\} d u} \rightarrow \Phi\left(z+\epsilon_{1}\right) e^{8 \epsilon_{1}}
\end{aligned}
$$

Similarly we prove that for all $\epsilon_{1}, N_{n} / D_{n} \geq \Phi\left(z-\epsilon_{1}\right) e^{-8 \epsilon_{1}}$, when $n$ is large enough, which terminates the proof of Theorem 2.1.
5. Proof of Theorem 2.2. Let $\beta>5 / 2$ and $\theta_{o, j}=c_{0} j^{-\left(\beta+\frac{1}{2}\right)}(\log j)^{-1}$. When the constant $c_{0}$ is chosen small enough, $\theta_{o} \in \Theta\left(\beta, L_{o}\right)$. In view of Corollary 3.1, the posterior mass on the events $\left\{(d, k, \theta):\left\|\theta-\bar{\theta}_{d, k}\right\| \geq 2 l_{0} \epsilon_{n}\right\}$ and $\left\{(d, k, \theta):\left|d-d_{o}\right| \geq w_{n}\right\}$ tends to zero in probability, and may be neglected. Moreover Lemma 3.1 implies that with posterior probability going to $1,\left\|\theta-\theta_{0}\right\| \lesssim(n / \log n)^{-(\beta-1 / 2) /(2 \beta+1)}$. However, within the $(k+1)$ dimensional FEXP-model, $\left\|\theta-\theta_{o}\right\|$ is minimized by setting $\theta_{j}=\theta_{o, j}(j=$ $0, \ldots, k)$, and for this choice of $\theta$ we have

$$
\left\|\theta-\theta_{o}\right\|^{2}=\sum_{l>k} \theta_{o, l}^{2} \gtrsim k^{-2 \beta}(\log k)^{-2}
$$

Consequently, the fact that $\left\|\theta-\theta_{0}\right\| \lesssim(n / \log n)^{-(\beta-1 / 2) /(2 \beta+1)}$ implies that $k>k_{n}^{\prime \prime}:=k_{l}(n / \log n)^{(\beta-1 / 2) /(\beta(2 \beta+1))}(\log n)^{-1 / \beta}$, for some constant $k_{l}$. We conclude that

$$
\Pi\left(k \leq k_{n}^{\prime \prime} \mid X\right)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1),
$$

and we can restrict our attention to $k>k_{n}^{\prime \prime}$.
We decompose $\Pi_{d}\left(\left|d-d_{o}\right| \leq k_{v} w_{n}(\log n)^{-1}, k>k_{n}^{\prime \prime} \mid X\right)$ as

$$
\begin{aligned}
& \sum_{m>k_{n}^{\prime \prime}} \Pi\left(\left|d-d_{o}\right| \leq k_{v} w_{n}(\log n)^{-1}, k=m \mid X\right) \\
& \quad=\sum_{m>k_{n}^{\prime \prime}} \Pi(k=m \mid X) \Pi_{m}\left(\left|d-d_{o}\right| \leq k_{v} w_{n}(\log n)^{-1} \mid X\right),
\end{aligned}
$$

where $\Pi_{m}\left(\left|d-d_{o}\right| \leq k_{v} w_{n}(\log n)^{-1} \mid X\right)$ is the posterior for $d$ within the FEXP-model of dimension $m+1$, i.e. $\Pi_{m}\left(\left|d-d_{o}\right| \leq k_{v} w_{n}(\log n)^{-1} \mid X\right):=$ $\Pi\left(\left|d-d_{o}\right| \leq k_{v} w_{n}(\log n)^{-1} \mid k=m, X\right)$.

To prove Theorem 2.2 it now suffices to show that

$$
\begin{align*}
& \sum_{k_{n}^{\prime \prime} \leq m \leq k_{n}^{\prime}} \Pi(k=m \mid X)=\Pi\left(k_{n}^{\prime \prime} \leq k \leq k_{n}^{\prime} \mid X\right) \xrightarrow{P_{\rho}} 1,  \tag{5.1}\\
& E_{0}^{n} \Pi_{k}\left(\left|d-d_{o}\right| \leq k_{v} w_{n}(\log n)^{-1} \mid X\right) \xrightarrow{P_{P}} 0, \quad \forall k_{n}^{\prime \prime} \leq k \leq k_{n}^{\prime} . \tag{5.2}
\end{align*}
$$

The convergence in (5.1) is a by-product of Theorem 1 in Rousseau and Kruijer (2011). In the remainder we prove (5.2). For every $k \leq k_{n}^{\prime}$ we can write, using the notation of (4.1),

$$
\begin{align*}
& \Pi_{k}\left(\left|d-d_{o}\right|<k_{v} w_{n}(\log n)^{-1} \mid X\right) \leq \frac{N_{n, k}}{D_{n, k}}  \tag{5.3}\\
& :=\frac{\int_{\left|d-d_{o}\right|<k_{v} w_{n}(\log n)^{-1}} \exp \left\{l_{n}(d, k)-l_{n}\left(d_{o}, k\right)+\log I_{n}(d, k)\right\} d \pi_{d}(d)}{\int_{\left|d-d_{o}\right|<w_{n}} \exp \left\{l_{n}(d, k)-l_{n}\left(d_{o}, k\right)+\log I_{n}(d, k)\right\} d \pi_{d}(d)}
\end{align*}
$$

Let $\delta_{2}>0$ and $A_{n}$ be the set of $X \in \mathbb{R}^{n}$ such that

$$
\left.\begin{array}{r}
\left|\log I_{n}(d, k)-\log I_{n}\left(d_{o}, k\right)\right| \leq \epsilon_{1}, \\
\left|l_{n}^{(1)}\left(d_{o}, k\right)-\mathcal{D}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)\right| \leq n^{\frac{1}{2}} k^{-\frac{1}{2}} \sqrt{\log n}, \\
\left|l_{n}^{(2)}(d, k)-\mathcal{D}\left(l_{n}^{(2)}\left(d_{o}, k\right)\right)\right| \leq \epsilon_{1} n^{-\left(2+\delta_{2}\right) /(2 \beta+1)}
\end{array}\right\}
$$

for all $\left|d-d_{o}\right| \leq w_{n}$ and $k_{n}^{\prime \prime} \leq k \leq k_{n}^{\prime}$. Compared to the definition of $A_{n}$ in the proof of Theorem 2.1, the constraints on $l_{n}^{(2)}(d, k)$ and $I_{n}$ are different. For the latter, recall from Lemma 3.5 that $\log I_{n}(d, k)=\log I_{n}\left(d_{o}, k\right)+\mathbf{o}_{\mathbf{P}}^{\mathbf{o}} \mathbf{( 1 )}$, uniformly over $d \in\left(d_{o}-w_{n}, d_{o}+w_{n}\right)$. As in the proof of Theorem 2.1, it now follows from Lemmas 3.4 and 3.5 that $P_{o}^{n}\left(A_{n}^{c}\right) \rightarrow 0$. We can write

$$
E_{0}^{n}\left[\frac{N_{n, k}}{D_{n, k}}\right] \leq P_{o}^{n}\left(A_{n}^{c}\right)+E_{0}^{n}\left[\frac{N_{n, k}}{D_{n, k}} 1_{A_{n}}\right],
$$

and bound $N_{n, k} / D_{n, k}$ pointwise for $X \in A_{n}$. Since when $k \in\left(k_{n}^{\prime \prime}, k_{n}^{\prime}\right)$,

$$
\frac{\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)^{2}}{2\left|l_{n}^{(2)}\left(d_{o}, k\right)\right|} n^{-\left(2+\delta_{2}\right) /(2 \beta+1)}=o(1)
$$

on $A_{n}$, for all $\delta_{2}>0$, analogous to (4.3) and (4.4), we find that for all $X \in A_{n}$, by definition of $b_{n}\left(d_{o}\right)$,
$l_{n}(d, k)-l_{n}\left(d_{o}, k\right)+\log I_{n}(d, k) \leq 2 \epsilon_{1}-\frac{\left|l_{n}^{(2)}\left(d_{o}, k\right)\right|}{2}\left(d-d_{o}-b_{n}\left(d_{o}\right)\right)^{2}+\frac{\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)^{2}}{2\left|l_{n}^{(2)}\left(d_{o}, k\right)\right|}$
$l_{n}(d, k)-l_{n}\left(d_{o}, k\right)+\log I_{n}(d, k) \geq-2 \epsilon_{1}-\frac{\left|l_{n}^{(2)}\left(d_{o}, k\right)\right|}{2}\left(d-d_{o}-b_{n}\left(d_{o}\right)\right)^{2}+\frac{\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)^{2}}{2\left|l_{n}^{(2)}\left(d_{o}, k\right)\right|}$,
when $n$ is large enough since $k>k_{n}^{\prime \prime}$. We now lower-bound $b_{n}\left(d_{o}\right)$ by bounding the terms on the right in (3.6) in Lemma 3.4. By construction of $\theta_{o}$ it follows that

$$
r_{k}^{-1} \sum_{j>k} j^{-1} \theta_{o, j}=c_{0} r_{k}^{-1} \sum_{j>k} j^{-\beta-\frac{3}{2}} /(\log j) \geq c k^{-\beta+\frac{1}{2}}(\log k)^{-1},
$$

for some $c>0$. Since $X \in A_{n}, 2 \mathcal{S}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right) /\left(n r_{k}\right) \leq 2 \sqrt{k / n} \sqrt{\log n}$. Since $k \leq k_{n}^{\prime}$, this bound is $o\left(k^{-\beta+\frac{1}{2}}(\log k)^{-1}\right)$. The last term in (3.6) is $o\left(n^{\epsilon-1}\right)$ when $\beta>5 / 2$, and hence this term is also $o\left(k^{-\beta-\frac{1}{2}}(\log k)^{-1}\right)$. Therefore, the last two terms in (3.6) are negligible with respect to $r_{k}^{-1} \sum_{j>k} j^{-1} \theta_{o, j}$. We deduce that $b_{n}\left(d_{o}\right) \geq c k^{-\beta+\frac{1}{2}}(\log k)^{-1} \geq c n^{-(2 \beta-1) /(4 \beta+2)}(\log n)^{-(2 \beta+3) /(4 \beta+2)}$ for $n$ large enough.

Consequently, when the constant $k_{v}$ is chosen sufficiently small, $\sqrt{n r_{k_{n}^{\prime}}}\left(b_{n}\left(d_{o}\right)-\right.$ $\left.k_{v} w_{n}(\log n)^{-1}\right) \geq\left(c-k_{v}\right) n^{1 /(4 \beta+2)}(\log n)^{-(\beta+1) /(2 \beta+1)}:=z_{n} \rightarrow \infty$. We now substitute the above bounds on $l_{n}(d, k)-l_{n}\left(d_{o}, k\right)+\log I_{n}(d, k)$ in the right hand side of (5.3), make the change of variables $u=d-d_{o}-b_{n}\left(d_{o}\right)$ and obtain

$$
\begin{aligned}
\frac{N_{n, k}}{D_{n, k}} & \leq e^{5 \epsilon_{1}} \frac{\int_{u \leq-k_{v} w_{n}(\log n)^{-1}-b_{n}\left(d_{o}\right)} \exp \left\{-\frac{n r_{k} u^{2}}{4}\right\} d u}{\int_{|u|<w_{n} / 2} \exp \left\{-\frac{n r_{k} u^{2}}{4}\right\} d u} \\
& \leq e^{5 \epsilon_{1}} \frac{\int_{v>z_{n}} \exp \left\{-\frac{v^{2}}{2}\right\} d v}{\int_{|v|<w_{n} \sqrt{n r_{k} / 8}} \exp \left\{-\frac{v^{2}}{2}\right\} d v}=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)
\end{aligned}
$$

This achieves the proof of Theorem 2.2.
6. Conclusion. In this paper we have derived conditions leading to a BVM type of result for the long memory parameter $d \in\left(-\frac{1}{2}, \frac{1}{2}\right)$ of a stationary Gaussian process, for the class of FEXP-priors. To our knowledge such a result has not been obtained before. The result implies in particular that asymptotically credible intervals for $d$ have good frequentist coverage.

A by-product of our results is that the most natural prior (Prior C) from a Bayesian perspective, which is also the prior leading to adaptive minimax rates under the loss function $l$ on $f$, leads to sub-optimal estimators in terms of $d$. Prior A leads to optimal estimators for $d$ however it is not adaptive. An interesting direction for future work would be to define an adaptiveminimax estimation procedure for $d$.

More broadly speaking, the approach considered here to derive the asymptotic posterior distribution of a finite dimensional parameter of interest in a semi-parametric problems could be used in other non - regular models, hence completing (not exhaustively) the recent works of Castillo (2010) and Bickel and Kleijn (2010).
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## APPENDIX A: PROOF OF LEMMA 3.4

We decompose the first derivative of $l_{n}(d, k)$ as $l_{n}^{(1)}(d, k)=\mathcal{S}\left(l_{n}^{(1)}(d, k)\right)+$ $\mathcal{D}\left(l_{n}^{(1)}(d, k)\right), \mathcal{S}\left(l_{n}^{(1)}(d, k)\right)$ being a centered quadratic form and $\mathcal{D}\left(l_{n}^{(1)}(d, k)\right)$ the remaining deterministic term. To simplify notations, in this proof we write $\mathcal{S}=\mathcal{S}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)$ and $\mathcal{D}=\mathcal{D}\left(l_{n}^{(1)}\left(d_{o}, k\right)\right)$. Using (1.6) (supplement) and defining $A=T_{n}^{-1}\left(f_{d_{o}, k}\right) T_{n}\left(H_{k} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)$, we find that

$$
\mathcal{D}=-\frac{1}{2} \operatorname{tr}\left[\left(T_{n}\left(f_{d_{o}, k}\right)-T_{n}\left(f_{o}\right)\right) A\right], \quad \mathcal{S}=\frac{1}{2}\left(X^{t} A X-\operatorname{tr}\left[T_{n}\left(f_{o}\right) A\right]\right) .
$$

From (1.4) and (1.8) in the supplement it follows that

$$
\begin{align*}
f_{o}-f_{d_{o}, k} & =f_{d_{o}, k}\left(e^{\Delta_{d_{o}, k}}-1\right)=\left(\Delta_{d_{o}, k}+\frac{1}{2} e^{\xi} \Delta_{d_{o}, k}^{2}\right) f_{d_{o}, k}  \tag{A.1}\\
& \left.=f_{d_{o}, k} O\left(k^{-\beta+1 / 2}\right), \quad \xi \in\left(0,\left(\Delta_{d_{o}, k}\right)_{+}\right)\right) .
\end{align*}
$$

Consequently, we have

$$
\begin{aligned}
\mathcal{D} & =\frac{1}{2} \operatorname{tr}\left[T_{n}\left(f_{d_{o}, k}\left(\Delta_{d_{o}, k}+O\left(\Delta_{d_{o}, k}^{2}\right)\right)\right) T_{n}^{-1}\left(f_{d_{o}, k}\right) T_{n}\left(H_{k} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)\right] \\
& =\frac{n}{4 \pi} \int_{-\pi}^{\pi} H_{k}(x)\left(\Delta_{d_{o}, k}(x)+O\left(\Delta_{d_{o}, k}^{2}(x)\right)\right) d x+\text { error } \\
& =\frac{n}{2} \sum_{j=k+1}^{\infty} \eta_{j} \theta_{o, j}+O\left(n k^{-2 \beta-1}\right)+\text { error. }
\end{aligned}
$$

The last equality follows from (1.9) and (1.11) in the supplement. We bound the error term using Lemma 2.4 (supplement) applied to $H_{k} f_{d_{o}, k}$ and $f_{d_{o}, k}$, whose Lipschitz constants are bounded by $O(k)$ and $O\left(k^{(3 / 2-\beta)_{+}}\right.$, respectively (see Lemma 3.1 in the supplement). Using that $\left\|\Delta_{d_{o}, k}\right\|_{\infty}=O\left(k^{-\beta+1 / 2}\right)$ (see (1.8) in the the supplement) we then find that the error is $O\left(k^{3 / 2-\beta} n^{\epsilon}\right)$.

The term $\mathcal{S}$ is a centered quadratic form with variance $\frac{1}{2}\left|T_{n}^{\frac{1}{2}}\left(f_{o}\right) A T_{n}^{\frac{1}{2}}\left(f_{o}\right)\right|^{2}$. Applying once more (A.1), we find that

$$
\begin{aligned}
\operatorname{tr}\left[\left(T_{n}\left(f_{o}\right) A\right)^{2}\right] & =\operatorname{tr}\left[\left(T_{n}^{-1}\left(f_{d_{o}, k}\right) T_{n}\left(H_{k} f_{d_{o}, k}\right)\right)^{2}\right]\left(1+O\left(\left\|\Delta_{d_{o}, k}\right\|_{\infty}\right)\right) \\
& =\frac{n}{2 \pi} \int_{-\pi}^{\pi} H_{k}^{2}(x) d x+O\left(n^{\epsilon} k+\left\|\Delta_{d_{o}, k}\right\|_{\infty}\right)=n r_{k}\left(1+o\left(n^{-\delta}\right)\right)
\end{aligned}
$$

where the term $n^{\epsilon} k$ comes from Lemma 2.4 in the supplement, associated to $f_{d_{o}, k}$ and $f_{d_{o}, k} H_{k}$. This proves the first equality in Lemma 3.4.

Similar to the decomposition of $l_{n}^{(1)}\left(d_{o}, k\right)$, we decompose the second derivative as $l_{n}^{(2)}(d, k)=\mathcal{D}\left(l_{n}^{(2)}(d, k)\right)-2 \mathcal{S}_{1}\left(l_{n}^{(2)}(d, k)\right)+\mathcal{S}_{2}\left(l_{n}^{(2)}(d, k)\right)$, where

$$
\mathcal{S}_{1}(d)=X^{t} A_{1, d} X-\operatorname{tr}\left[T_{n}\left(f_{o}\right) A_{1, d}\right], \quad \mathcal{S}_{2}(d)=X^{t} A_{2, d} X-\operatorname{tr}\left[T_{n}\left(f_{o}\right) A_{2, d}\right],
$$

$$
\begin{align*}
\mathcal{D}_{2}(d): & =\mathcal{D}\left(l_{n}^{(2)}(d, k)\right)  \tag{A.2}\\
= & -\frac{1}{2} \operatorname{tr}\left[T_{n}\left(f_{d, k}\right) A_{1, d}\right]+\operatorname{tr}\left[\left(T_{n}\left(f_{d, k}\right)-T_{n}\left(f_{d_{o}, k}\right)\right)\left(A_{1, d}-\frac{1}{2} A_{2, d}\right)\right] \\
& +\operatorname{tr}\left[\left(T_{n}\left(f_{d_{o}, k}\right)-T_{n}\left(f_{o}\right)\right)\left(A_{1, d}-\frac{1}{2} A_{2, d}\right)\right] \\
A_{1, d}= & T_{n}^{-1}\left(f_{d, k}\right)\left(T_{n}\left(H_{k} f_{d, k}\right) T_{n}^{-1}\left(f_{d, k}\right)\right)^{2}, A_{2, d}=T_{n}^{-1}\left(f_{d, k}\right) T_{n}\left(H_{k}^{2} f_{d, k}\right) T_{n}^{-1}\left(f_{d, k}\right) .
\end{align*}
$$

To control $\mathcal{D}_{2}(d)$ we use a first order Taylor expansion around $d_{o}$, implying that $\mathcal{D}_{2}(d)=\mathcal{D}_{2}\left(d_{o}\right)+O\left(\mid d-d_{o}\right) \sup _{\left|d^{\prime}-d_{o}\right| \leq \bar{v}_{n}}\left|\mathcal{D}_{2}^{\prime}\left(d^{\prime}\right)\right|$. First we study $\mathcal{D}_{2}\left(d_{o}\right)$. At $d=d_{o}$, the right-hand side of (A.2) equals

$$
\begin{gather*}
-\frac{n}{4 \pi} \int_{-\pi}^{\pi} H_{k}^{2}(x)\left(1+\left(e^{\Delta_{d_{o}, k}}-1\right)\right) d x+O\left(k n^{\epsilon}\right)  \tag{A.3}\\
\quad=-\frac{-n r_{k}}{2}\left(1+O\left(k^{-\beta+1 / 2}+k^{2} / n^{1-\epsilon}\right)\right)
\end{gather*}
$$

The $O\left(k n^{\epsilon}\right)$ term is obtained from Lemma 2.4 (supplement), applied to $f_{2 j}=H_{k} f_{d_{o}, k}$ and $f_{2 j-1}=f_{d_{o}, k}$, with Lipschitz constants $O(k)$ for the former and $O\left(k^{(3 / 2-\beta)_{+}}\right)$for the latter, together with the bound $\left\|\Delta_{d_{o}, k}\right\|_{\infty}=$ $O\left(k^{-\beta+1 / 2}\right)$. Using

$$
\begin{align*}
A_{1, d}^{\prime}= & -3\left(T_{n}^{-1}\left(f_{d, k}\right) T_{n}\left(H_{k} f_{d, k}\right)\right)^{3} T_{n}^{-1}\left(f_{d, k}\right)  \tag{A.4}\\
& +2 T_{n}^{-1}\left(f_{d, k}\right) T_{n}\left(H_{k}^{2} f_{d, k}\right) T_{n}^{-1}\left(f_{d, k}\right) T_{n}\left(H_{k} f_{d, k}\right) T_{n}^{-1}\left(f_{d, k}\right)
\end{align*}
$$

and a similar expression for the derivative of $d \rightarrow A_{2, d}$, it follows that

$$
\begin{aligned}
& \left|\mathcal{D}_{2}^{\prime}\left(d^{\prime}\right)\right| \lesssim \operatorname{tr}\left[\left(T_{n}\left(f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d^{\prime}, k}\right)+I_{n}\right)\left(T_{n}\left(\left|H_{k}\right| f_{d^{\prime}, k}\right) T_{n}^{-1}\left(f_{d^{\prime}, k}\right)\right)^{3}\right] \\
& \quad+\operatorname{tr}\left[\left(T_{n}\left(f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d^{\prime}, k}\right)+I_{n}\right)\left(T_{n}\left(\left|H_{k}\right| f_{d^{\prime}, k}\right) T_{n}^{-1}\left(f_{d^{\prime}, k}\right)\right) T_{n}\left(H_{k}^{2} f_{d^{\prime}, k}\right) T_{n}^{-1}\left(f_{d^{\prime}, k}\right)\right] \\
& \quad+\operatorname{tr}\left[\left(T_{n}\left(f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d^{\prime}, k}\right)+I_{n}\right) T_{n}\left(\left|H_{k}\right|^{3} f_{d^{\prime}, k}\right) T_{n}^{-1}\left(f_{d^{\prime}, k}\right)\right]
\end{aligned}
$$

We control the first term of the right hand side of the above inequality, the second and third terms are controlled similarly. Note first that

$$
\begin{align*}
& \operatorname{tr}\left[T_{n}\left(f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d^{\prime}, k}\right)\left(T_{n}\left(\left|H_{k}\right| f_{d^{\prime}, k}\right) T_{n}^{-1}\left(f_{d^{\prime}, k}\right)\right)^{3}\right]  \tag{A.5}\\
& =\left|T_{n}^{\frac{1}{2}}\left(f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d^{\prime}, k}\right) T_{n}\left(\left|H_{k}\right| f_{d^{\prime}, k}\right) T_{n}^{-1}\left(f_{d^{\prime}, k}\right) T_{n}^{\frac{1}{2}}\left(\left|H_{k}\right| f_{d^{\prime}, k}\right)\right|^{2} \\
& \leq\left\|T_{n}^{\frac{1}{2}}\left(f_{d_{o}, k}\right) T_{n}^{-\frac{1}{2}}\left(f_{d^{\prime}, k}\right)\right\|^{2} \\
& \quad \times\left\|T_{n}^{-\frac{1}{2}}\left(f_{d^{\prime}, k}\right) T_{n}^{\frac{1}{2}}\left(\left|H_{k}\right| f_{d^{\prime}, k}\right)\right\|^{2}\left|T_{n}^{-\frac{1}{2}}\left(f_{d^{\prime}, k}\right) T_{n}\left(\left|H_{k}\right| f_{d^{\prime}, k}\right) T_{n}^{-\frac{1}{2}}\left(f_{d^{\prime}, k}\right)\right|^{2} \\
& \lesssim n^{\epsilon}\left|T_{n}^{-\frac{1}{2}}\left(f_{d^{\prime}, k}\right) T_{n}\left(\left|H_{k}\right| f_{d^{\prime}, k}\right) T_{n}^{-\frac{1}{2}}\left(f_{d^{\prime}, k}\right)\right|^{2},
\end{align*}
$$

where the last inequality comes from Lemma 2.3 in the supplement. Note also that
$|x|^{-2 d^{\prime}} \lesssim f_{d^{\prime}}(x) \lesssim|x|^{-2 d^{\prime}} \quad$ and $T_{n}\left(\left|H_{k}\right| f_{d^{\prime}, k}\right) \lesssim T_{n}\left(\left|H_{k}\right||x|^{-2 d^{\prime}}\right), \quad T_{n}\left(\mid f_{d^{\prime}, k}\right) \gtrsim T_{n}\left(|x|^{-2 d^{\prime}}\right)$
and replace $f_{d^{\prime}}$ by $|x|^{-2 d^{\prime}}$ in (A.5), then

$$
\left|T_{n}^{-1 / 2}\left(f_{d^{\prime}, k}\right) T_{n}\left(\left|H_{k}\right| f_{d^{\prime}, k}\right) T_{n}^{-1 / 2}\left(f_{d^{\prime}, k}\right)\right|^{2} \lesssim\left(\frac{n}{k}+O(k)\right)
$$

using Lemma 2.4 in the supplement associated to $\left|H_{k}\right||x|^{-2 d^{\prime}}$ which has Lipschitz constant $k$. This leads to $\mathcal{D}_{2}\left(d^{\prime}\right)=O\left(n^{\epsilon} \frac{n}{k}\right)$, which implies that for all $\beta>1$,

$$
\mathcal{D}_{2}(d)=\mathcal{D}_{2}\left(d_{o}\right)+\mathbf{o}\left(\left|d-d_{o}\right| n^{\epsilon+1} k^{-1}\right)=-\frac{n r_{k}}{2}\left(1+\mathbf{o}\left(n^{-\delta}\right)\right) .
$$

For the stochastic terms in $l_{n}^{(2)}(d, k)$ we need a chaining argument to control the supremum over $d \in\left(d_{o}-\bar{v}_{n}, d_{o}+\bar{v}_{n}\right)$. We show that for all $\epsilon^{\prime}>0$ and $\gamma_{n}=n^{\frac{1}{2}+\epsilon^{\prime}} k^{-\frac{1}{2}}$,

$$
\begin{equation*}
P_{o}^{n}\left(\sup _{\left|d-d_{o}\right| \leq \bar{v}_{n}}\left|\mathcal{S}_{1}(d)\right|>\gamma_{n}\right)=o(1) \tag{A.6}
\end{equation*}
$$

i.e. that $\mathcal{S}_{1}(d)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(\gamma_{n}\right)$. The same can be shown for $\mathcal{S}_{2}(d)$ using exactly the same arguments. Consider a covering of $\left(d_{o}-\bar{v}_{n}, d_{o}+\bar{v}_{n}\right)$ by balls of radius $n^{-1}$ centered at $d_{j}, j=1, \ldots, J_{n}$ with $J_{n} \leq 2 \bar{v}_{n} n$. Then

$$
\sup _{\left|d-d_{o}\right|<\bar{v}_{n}}\left|\mathcal{S}_{1}(d)\right| \leq \max _{j}\left|\mathcal{S}_{1}\left(d_{j}\right)\right|+\sup _{\left|d-d^{\prime}\right| \leq n^{-1}}\left|\mathcal{S}_{1}(d)-\mathcal{S}_{1}(d)\right|,
$$

and

$$
\begin{align*}
P_{o}^{n}\left(\sup _{\left|d-d_{o}\right| \leq \bar{v}_{n}}\left|\mathcal{S}_{1}(d)\right|>\gamma_{n}\right) \leq & P_{o}^{n}\left(\sup _{\left|d-d^{\prime}\right| \leq n^{-1}}\left|\mathcal{S}_{1}(d)-\mathcal{S}_{1}\left(d^{\prime}\right)\right|>\frac{1}{2} \gamma_{n}\right)  \tag{A.7}\\
& +J_{n} \max _{1 \leq j \leq J_{n}} P_{o}^{n}\left(\left|\mathcal{S}_{1}\left(d_{j}\right)\right|>\frac{1}{2} \gamma_{n}\right)
\end{align*}
$$

To control the first term on the right in (A.7), note that for a standard normal vector $Z$ and some $d^{*} \in\left(d, d^{\prime}\right)$,

$$
\mathcal{S}_{1}(d)-\mathcal{S}_{1}\left(d^{\prime}\right)=\left(d-d^{\prime}\right)\left(Z^{t} T_{n}^{\frac{1}{2}}\left(f_{o}\right) A_{1, d^{*}}^{\prime} T_{n}^{\frac{1}{2}}\left(f_{o}\right) Z-\operatorname{tr}\left[T_{n}\left(f_{o}\right) A_{1, d^{*}}^{\prime}\right]\right),
$$

with $A_{1, d}^{\prime}$ as in (A.4). Using Lemma 2.3 (supplement) and the fact that $\|A B\| \leq\|A\|\|B\|$ for all matrices $A$ and $B$, it follows that $\left\|T_{n}^{\frac{1}{2}}\left(f_{o}\right) A_{1, d^{*}}^{\prime} T_{n}^{\frac{1}{2}}\left(f_{o}\right)\right\|=$ $O\left(n^{\epsilon}\right)$, and hence $Z^{t} T_{n}^{\frac{1}{2}}\left(f_{o}\right) A_{1, d^{*}}^{\prime} T_{n}^{\frac{1}{2}}\left(f_{o}\right) Z \leq Z^{t} Z\left\|T_{n}^{\frac{1}{2}}\left(f_{o}\right) A_{1, d^{*}}^{\prime} T_{n}^{\frac{1}{2}}\left(f_{o}\right)\right\|=$ $O\left(n^{\epsilon}\right) Z^{t} Z$. Similarly, it follows that $\operatorname{tr}\left[T_{n}\left(f_{o}\right) A_{1, d^{*}}^{\prime}\right] \lesssim n$. Consequently, when $\epsilon=\epsilon^{\prime} / 2$ we have

$$
\left|\mathcal{S}_{1}(d)-\mathcal{S}_{1}\left(d^{\prime}\right)\right| \lesssim n^{-1}\left(Z^{t} Z n^{\epsilon}+n\right),
$$

uniformly over all $d, d^{\prime}$ such that $\left|d-d^{\prime}\right| \leq n^{-1}$. Since $1=o\left(\gamma_{n}\right)$,

$$
P_{o}^{n}\left(\sup _{\left|d-d^{\prime}\right| \leq n^{-1}}\left|\mathcal{S}_{1}(d)-\mathcal{S}_{1}\left(d^{\prime}\right)\right| \geq \frac{1}{2} \gamma_{n}\right) \leq P\left(Z^{t} Z>n^{1-\epsilon} \gamma_{n} / 4\right)=o(1)
$$

To bound the last term in (A.7), we apply Lemma 1.3 (supplement) to $\left(Z^{t} A Z-\operatorname{tr}[A]\right)|A|^{-1}$, with $A=T_{n}^{\frac{1}{2}}\left(f_{o}\right) A_{1, d} T_{n}^{\frac{1}{2}}\left(f_{o}\right)$ since as seen previously $|A|^{2}=\mathbf{O}(n / k)=\mathbf{o}\left(\gamma_{n}^{2} n^{-2 \alpha}\right)$ for $\alpha$ small enough, it follows that

$$
P_{o}^{n}\left(\mathcal{S}_{1}(d) \geq \frac{1}{2} \gamma_{n}\right) \leq e^{-n^{\alpha} / 8}
$$

Since $J_{n}$ increases only polynomially with $n$, this finishes the proof of (A.6).

## APPENDIX B: CONTROL OF THE DERIVATIVES IN $\theta$ ON THE LOG-LIKELIHOOD

Before stating Lemma B. 1 we first give a general expression for the derivatives of $l_{n}(d, k, \theta)$ with respect to $\theta$. For all $j \geq 1$ and $l=\left(l_{1}, \ldots, l_{j}\right) \in$ $\{0,1, \ldots, k\}^{j}$, let $\sigma=(\sigma(1), \ldots, \sigma(|\sigma|))$ be a partition of $\{1, \ldots, j\}$. Let $|\sigma|$ be the number of subsets in this partition and $\sigma(i)$ the $i$ th subset of $\{1, \ldots, j\}$ in the partition $\sigma$. Denoting $l_{\sigma(i)}$ the vector $\left(l_{t}, t \in \sigma(i)\right)$, we can write

$$
\nabla_{l_{\sigma(i)}} f_{d, k, \theta}(x)=\prod_{t \in \sigma(i)} \cos \left(l_{t} x\right) f_{d, k, \theta}(x) .
$$

For notational ease we write $\nabla_{\sigma(i)} f_{d, k, \theta}:=\nabla_{l_{\sigma(i)}} f_{d, k, \theta}$. The derivative $\frac{\partial^{j} l_{n}(d, k, \theta)}{\partial \theta_{l_{1}} \ldots \partial \theta_{l_{j}}}$ can now be written in terms of the matrices

$$
\begin{equation*}
B_{\sigma}(d, \theta)=\prod_{i=1}^{|\sigma|} B_{\sigma(i)}(d, \theta), \quad B_{\sigma(i)}(d, \theta)=T_{n}\left(\nabla_{\sigma(i)} f_{d, k, \theta}\right) T_{n}^{-1}\left(f_{d, k, \theta}\right) \tag{B.1}
\end{equation*}
$$

There exist constants $b_{\sigma}, c_{\sigma}$ and $d_{\sigma}$ such that

$$
\begin{align*}
& \frac{\partial^{j} l_{n}(d, k, \theta)}{\partial \theta_{l_{1}} \ldots \partial \theta_{l_{j}}}  \tag{B.2}\\
& =\sum_{\sigma \in \mathcal{S}_{j}} b_{\sigma}\left(X^{t} T_{n}^{-1}\left(f_{d, k, \theta}\right) B_{\sigma}(d, \theta) X-\operatorname{tr}\left[T_{n}\left(f_{o}\right) T_{n}^{-1}\left(f_{d, k, \theta}\right) B_{\sigma}(d, \theta)\right]\right) \\
& +\sum_{\sigma \in \mathcal{S}_{j}} c_{\sigma} \operatorname{tr}\left[B_{\sigma}(d, \theta)\right]+\sum_{\sigma \in \mathcal{S}_{j}} d_{\sigma} \operatorname{tr}\left[\left(T_{n}\left(f_{o}\right) T_{n}^{-1}\left(f_{d, k, \theta}\right)-I_{n}\right) B_{\sigma}(d, \theta)\right],
\end{align*}
$$

where $\mathcal{S}_{j}$ is the set of partitions of $\{1, \ldots, j\}$. For the first two derivatives $(j=1,2)$ the values of the constants $b_{\sigma}, c_{\sigma}$ and $d_{\sigma}$ are given below in Lemmas B. 4 and B.5. For the higher order derivatives these values are not important for our purpose; we will only need that for any $j \geq 1$, the constant $c_{\sigma}$ is zero if $|\sigma|=1$.

The following lemma states that $l_{n}(d, k, \theta)-l_{n}(d, k)$ is the sum of a Taylorapproximation $\sum_{j=1}^{J} \frac{\left(\theta-\bar{\theta}_{d, k}\right)^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}$ and terms whose dependence on $d$ can be negligible. Since the proof is involved, some of the technical details are treated in Lemmas B. 2 and B.3.

Lemma B.1. Given $\beta>1$, let $k \leq k_{n}$ and let $d$ and $\theta$ be such that $l\left(f_{o}, f_{d, k, \theta}\right) \leq l_{0}^{2} \delta_{n}^{2}$. Then there exists an integer $J$ and a constant $\epsilon>0$ such
that uniformly over $d \in\left(d_{o}-\bar{v}_{n}, d_{o}+\bar{v}_{n}\right)$ and $\theta \in B_{k}\left(\bar{\theta}_{d, k}, 2 l_{0} \delta_{n}\right)$,

$$
\begin{align*}
l_{n}(d, k, \theta)-l_{n}(d, k)= & \sum_{j=1}^{J} \frac{\left(\theta-\bar{\theta}_{d, k}\right)^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!} \\
& +\left(d-d_{o}\right) \sum_{j=2}^{J} \frac{1}{j!} g_{n, j}\left(\theta-\bar{\theta}_{d, k}\right)+S_{n}(d), \tag{B.3}
\end{align*}
$$

where, for $u=\theta-\bar{\theta}_{d, k}$,
(B.4)

$$
\begin{gathered}
g_{n, j}(u)=\sum_{l_{1}, \ldots, l_{j}=0}^{k} u_{l_{1}} \cdots u_{l_{j}} \sum_{\sigma \in \mathcal{S}_{j}}\left(c_{\sigma} \operatorname{tr}\left[T_{1, \sigma}\left(d_{o}, k\right)\right]+d_{\sigma} \operatorname{tr}\left[T_{2, \sigma}\left(d_{o}, k\right)\right]\right), \\
T_{1, \sigma}\left(d_{o}, k\right)=\sum_{i=1}^{|\sigma|}\left(\prod_{l<i} T_{n}\left(\nabla_{\sigma(l)} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)\right) \times \\
{\left[T_{n}\left(\nabla_{\sigma(i)} f_{d_{o}, k} H_{k}\right)-T_{n}\left(H_{k} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right) T_{n}\left(\nabla_{\sigma(i)} f_{d_{o}, k}\right)\right] \times} \\
T_{n}^{-1}\left(f_{d_{o}, k}\right)\left(\prod_{l>i} T_{n}\left(\nabla_{\sigma(l)} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)\right), \\
T_{2, \sigma}\left(d_{o}, k\right)=-T_{n}\left(H_{k} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right) B_{\sigma}\left(d_{o}, \bar{\theta}_{d_{o}, k}\right),
\end{gathered}
$$

and $S_{n}(d)$ denotes any term of order

$$
\begin{equation*}
S_{n}(d)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(\frac{\left|d-d_{o}\right| n^{\frac{1}{2}-\delta}}{\sqrt{k}}\right)+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(\left(d-d_{o}\right)^{2} \frac{n^{1-\delta}}{k}\right) \tag{B.5}
\end{equation*}
$$

When $\beta>5 / 2$ and $k \leq k_{n}^{\prime}$, we can choose $J=2$, and (B.3) simplifies to

$$
\begin{equation*}
l_{n}(d, k, \theta)-l_{n}(d, k)=\sum_{j=1}^{2} \frac{\left(\theta-\bar{\theta}_{d, k}\right)^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1) \tag{B.6}
\end{equation*}
$$

Proof. Recall that by (3.7),

$$
\begin{align*}
& l_{n}(d, k, \theta)-l_{n}(d, k)=\sum_{j=1}^{J} \frac{\left(\theta-\bar{\theta}_{d, k}\right)^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!} \\
& \quad+\sum_{j=1}^{J} \frac{\left(\theta-\bar{\theta}_{d, k}\right)^{(j)} \nabla^{j}\left(l_{n}(d, k)-l_{n}\left(d_{o}, k\right)\right)}{j!}+R_{J+1, d}(\theta) . \tag{B.7}
\end{align*}
$$

To prove (B.3) we first show that, writing $u=\theta-\bar{\theta}_{d, k}$,

$$
\begin{align*}
\sum_{j=1}^{J} & \frac{u^{(j)} \nabla^{j}\left(l_{n}(d, k)-l_{n}\left(d_{o}, k\right)\right)}{j!}  \tag{B.8}\\
& =\sum_{j=1}^{J} \frac{1}{j!} \sum_{l_{1}, \ldots, l_{j}=0}^{k} u_{l_{1}} \ldots u_{l_{j}}\left(\frac{\partial^{j} l_{n}\left(d, k, \bar{\theta}_{d, k}\right)}{\partial \theta_{l_{1}} \ldots \partial \theta_{l_{j}}}-\frac{\partial^{j} l_{n}\left(d_{o}, k, \bar{\theta}_{d_{o}, k}\right)}{\partial \theta_{l_{1}} \ldots \partial \theta_{l_{j}}}\right) \\
& =\left(d-d_{o}\right) \sum_{j=1}^{J} \frac{1}{j!} g_{n, j}(u)+O\left(S_{n}(d)\right) .
\end{align*}
$$

This result is combined with (B.7) and Lemma B. 3 below, by which $g_{n, 1}(u)=$ $O\left(S_{n}(d)\right)$. It then follows that $l_{n}(d, k, \theta)-l_{n}(d, k)$ equals
$\sum_{j=1}^{J} \frac{\left(\theta-\bar{\theta}_{d, k}\right)^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}+\left(d-d_{o}\right) \sum_{j=2}^{J} \frac{1}{j!} g_{n, j}(u)+R_{J+1, d}(\theta)+O\left(S_{n}(d)\right)$.
The final step is to prove that $R_{J+1, d}(\theta)$ is $\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)$ and hence $O\left(S_{n}(d)\right)$; to this end $J$ needs to be sufficiently large.

First we prove (B.8). For the factors $\frac{\partial^{j}}{\partial \theta} l_{n}\left(d, k, \bar{\theta}_{d, k}\right)-\frac{\partial^{j}}{\partial \theta} l_{n}\left(d, k, \bar{\theta}_{d_{o}, k}\right)$ we substitute (B.2). In Lemma B. 2 below we give expressions for each of the terms therein, which we substitute in (B.8). The main terms are $\left(d-d_{o}\right) \operatorname{tr}\left[T_{1, \sigma}\left(d_{o}, k\right)\right]$ and $\left(d-d_{o}\right) \operatorname{tr}\left[T_{2, \sigma}\left(d_{o}, k\right)\right]$ in (B.13) and (B.14), which after substitution in (B.8) give the term $\left(d-d_{o}\right) \sum_{j=1}^{J} \frac{1}{j!} g_{n, j}(u)$ on the right. The other terms in (B.12)-(B.14) that enter (B.8) through (B.2) are $O\left(S_{n}(d)\right)$. This is due to the summation over $u_{l_{1}}, \ldots, u_{l_{j}}$ in (B.8), and the Cauchy-Schwarz inequality by which

$$
\begin{equation*}
\left|\sum_{l_{1}, \ldots, l_{j}=0}^{k} u_{l_{1}} \ldots u_{l_{j}}\right| \leq(\sqrt{k}\|u\|)^{j} \leq\left(2 l_{0} \sqrt{k} \delta_{n}\right)^{j}=o\left(n^{-\delta}\right), \tag{B.9}
\end{equation*}
$$

for some $\delta>0$, as $\|u\| \leq 2 l_{0} \delta_{n}$ and (B.8) is proved. We now control $R_{J+1, d}(\theta)$.

Combining (3.8) and the first inequality in (B.9), we obtain

$$
\left|R_{J+1, d}(\theta)\right| \leq \frac{1}{(J+1)!}\left(\sqrt{k} \delta_{n}\right)^{J+1} \max _{l_{1}, \ldots, l_{J+1}\left\|\tilde{\theta}-\bar{\theta}_{d, k}\right\| \leq 2 l_{0} \delta_{n}} \sup \left|\frac{\partial^{J+1} l_{n}(d, k, \tilde{\theta})}{\partial \theta_{l_{1}} \ldots \partial \theta_{l_{J+1}}}(x)\right| .
$$

We give a direct bound on this derivative using (B.2). For all partitions $\sigma$ of $\left\{l_{1}, \ldots, l_{J+1}\right\}$ and all $\left(l_{1}, \ldots, l_{J+1}\right) \in\{1, \ldots, k\}^{J+1}$, we bound $\left\|B_{\sigma(i)}(d, \theta)\right\|$,
using $\left\|B_{\sigma(i)}(d, \theta)\right\| \leq\left\|T_{n}^{\frac{1}{2}}\left(\nabla_{\sigma(i)} f_{d, k, \theta}\right) T_{n}^{-\frac{1}{2}}\left(f_{d, k, \theta}\right)\right\|^{2}$ (see (1.6)). We bound $\left\|T_{n}^{\frac{1}{2}}\left(\nabla_{\sigma(i)} f_{d, k, \theta}\right) T_{n}^{-\frac{1}{2}}\left(f_{d, k, \theta}\right)\right\|$ by application of Lemma 2.3 (supplement) with $f=f_{d, k, \theta}$ and $g=\nabla_{\sigma(i)} f_{d, k, \theta}$. The constant $M$ in this lemma is bounded by

$$
\sum_{j=0}^{k}\left|\theta_{j}\right| \leq \sum_{j=0}^{k}\left|\left(\bar{\theta}_{d, k}\right)_{j}\right|+\sum_{j=0}^{k}\left|\theta_{j}-\left(\bar{\theta}_{d, k}\right)_{j}\right| \leq 2 \sqrt{L}+\sqrt{k}\left\|\theta-\bar{\theta}_{d, k}\right\|=O(1)
$$

since $\sum_{i=0}^{k}\left|\left(\bar{\theta}_{d, k}\right)_{i}\right| \leq 2 \sqrt{L}$ (by Lemma 3.3) and $\left\|\theta-\bar{\theta}_{d, k}\right\| \leq \delta_{n}$. Consequently, Lemma 2.3 (supplement) implies that

$$
\begin{equation*}
\left\|B_{\sigma(i)}(d, \theta)\right\| \leq K \tag{B.10}
\end{equation*}
$$

where $K$ depends only on $L, L_{o}$ and not on $n, d$ nor $\theta$. From the relations in (1.6) and the definition of $B_{\sigma}$ it follows that for any $\sigma, d, \theta$,

$$
\begin{gathered}
\left|X^{t} B_{\sigma}(d, \theta) X\right| \leq X^{t} T_{n}^{-1}\left(f_{o}\right) X K^{|\sigma|}\left\|T_{n}^{\frac{1}{2}}\left(f_{o}\right) T_{n}^{-\frac{1}{2}}\left(f_{d, k}\right)\right\|^{2} \leq X^{t} T_{n}^{-1}\left(f_{o}\right) X K^{|\sigma|} n^{\epsilon}, \\
\left|\operatorname{tr}\left[B_{\sigma}(d, \theta)\right]\right| \leq n K^{|\sigma|}\left\|T_{n}^{\frac{1}{2}}\left(f_{o}\right) T_{n}^{-\frac{1}{2}}\left(f_{d, k}\right)\right\|^{2} \leq n^{1+\epsilon} K^{|\sigma|} .
\end{gathered}
$$

Therefore we have the bound

$$
\begin{equation*}
\left|R_{J+1}(d, \theta)\right| \leq C K^{J+1} n^{\epsilon}\left(\sqrt{k}\left\|\theta-\bar{\theta}_{d, k}\right\|\right)^{J+1}\left(X^{t} T_{n}^{-1}\left(f_{o}\right) X+n\right) . \tag{B.11}
\end{equation*}
$$

Since $k \leq k_{n},\left\|\theta-\bar{\theta}_{d, k}\right\| \leq \delta_{n}$ and the term $X^{t} T_{n}^{-1}\left(f_{o}\right) X$ in (B.11) is the sum of $n$ independent standard normal variables, there is a constant $c>0$ such that
provided we choose $J$ such that $(J+1)(1-1 / \beta)>2$. This concludes the proof of (B.3).

To prove (B.6) we first show that for $J=2,\left|R_{J+1}(d, \theta)\right|=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)$. Since $k \leq k_{n}^{\prime}, \beta>5 / 2$ and $\left\|\theta-\bar{\theta}_{d, k}\right\| \leq 2 l_{0} \epsilon_{n}$, we can choose $J+1=3>$ $(2 \beta+1) /\left(\beta-\frac{1}{2}\right)$, and the preceding inequality becomes

$$
P_{o}\left(\sup _{\left|d-d_{o}\right| \leq w_{n}\left\|\theta-\bar{\theta}_{d, k}\right\| \leq 2 l_{0} \epsilon_{n}} \sup \left|R_{3}(d, \theta)\right|>n^{-\epsilon}\right) \leq e^{-c n} .
$$

Combining this result with (B.8), it only remains to be shown that ( $d-$ $\left.d_{o}\right) g_{n, 1}(u)$ and $\left(d-d_{o}\right) g_{n, 2}(u)$ are $\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)$. Recall from Corollary 3.1 that $\left|d-d_{o}\right|=o\left(n^{\epsilon-(\beta-1 / 2) /(2 \beta+1)}\right)$ for all $\epsilon>0$. Consequently,
$\left|d-d_{o}\right| n^{\frac{1}{2}} k^{-\frac{1}{2}}=o\left(n^{\epsilon+1 /(2 \beta+1)} k^{-\frac{1}{2}}\right)=o\left(\left(\sqrt{k} \epsilon_{n}\right)^{-1}\right), \quad\left(d-d_{o}\right)^{2} \frac{n}{k}=o\left(\left(\sqrt{k} \epsilon_{n}\right)^{-1}\right)$
for all $\beta>2$. This implies that $S_{n}(d)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)$ and that, by Lemma B.3, $\left(d-d_{o}\right) g_{n, 1}(u)=\mathbf{o}(1)$. Also, for all $\left(l_{1}, l_{2}\right) \in\{1, \ldots, k\}^{j}$ and all partitions $\sigma$ of $\left(l_{1}, l_{2}\right)$, the limiting integral of $\operatorname{tr}\left[T_{1, \sigma}\right]$ is equal to 0 . Since $\beta>5 / 2$ the Lipschitz constants of the functions $f_{d_{o}, k}$ or $f_{o}$ are $O(1)$, so that Lemma 2.4 (supplement) implies $\operatorname{tr}\left[T_{1, \sigma}\left(d_{o}, k\right)\right]=O\left(n^{\epsilon} k\right)$. Similarly,

$$
\operatorname{tr}\left[T_{2, \sigma}\left(d_{o}, k\right)\right]=\frac{n}{2 \pi} \int_{-\pi}^{\pi} H_{k}(x) \cos \left(l_{1} x\right) \cos \left(l_{j} x\right) d x+O\left(n^{\epsilon} k\right) .
$$

Thus we have
$\left(d-d_{o}\right) g_{n, 2}(u)=\frac{n\left(d-d_{o}\right)}{2 \pi} \sum_{l_{1}, l_{2}=0}^{k} u_{l_{1}} u_{l_{2}} \int_{-\pi}^{\pi} H_{k}(x) \cos \left(l_{1} x\right) \cos \left(l_{j} x\right) d x+\mathbf{o}(1)$,
which is $\mathbf{o}(1)$. This completes the proof of Lemma B.1.
The proof of the following lemma is given in section 4 of the supplement.
Lemma B.2. Let $W_{\sigma}(d)$ denote any of the quadratic forms

$$
X^{t} T_{n}^{-1}\left(f_{d, k}\right) B_{\sigma}\left(d, \bar{\theta}_{d, k}\right) X-\operatorname{tr}\left[T_{n}\left(f_{o}\right) T_{n}^{-1}\left(f_{d, k}\right) B_{\sigma}\left(d, \bar{\theta}_{d, k}\right)\right]
$$

in (B.2). For any $j \leq J,\left(l_{1}, \ldots, l_{j}\right) \in\{0, \ldots, k\}^{j}$ and $\sigma \in \mathcal{S}_{j}$, we have

$$
\begin{align*}
& \left|W_{\sigma}(d)-W_{\sigma}\left(d_{o}\right)\right|=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(\left|d-d_{o}\right| n^{\frac{1}{2}+\epsilon} k^{-\frac{1}{2}}\right),  \tag{B.12}\\
\operatorname{tr} & {\left[B_{\sigma}\left(d, \bar{\theta}_{d, k}\right)\right]-\operatorname{tr}\left[B_{\sigma}\left(d_{o}, \bar{\theta}_{d_{o}}\right)\right] } \\
= & \left(d-d_{o}\right) \operatorname{tr}\left[T_{1, \sigma}\left(d_{o}, k\right)\right]+\left(d-d_{o}\right)^{2} \mathbf{o}\left(n^{\epsilon+\frac{1}{2}} k^{-\frac{1}{2}+(1-\beta / 2)_{+}}\right)  \tag{B.13}\\
= & \left(d-d_{o}\right) \operatorname{tr}\left[T_{1, \sigma}\left(d_{o}, k\right)\right]+\left(d-d_{o}\right)^{2} \mathbf{o}\left(n^{1-\delta} / k\right),
\end{align*}
$$

$$
\begin{align*}
& \operatorname{tr}\left[\left(T_{n}\left(f_{o}\right) T_{n}^{-1}\left(f_{d, k}\right)-I_{n}\right) B_{\sigma}\left(d, \bar{\theta}_{d, k}\right)\right]-\operatorname{tr}\left[\left(T_{n}\left(f_{o}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)-I_{n}\right) B_{\sigma}\left(d_{o}, \bar{\theta}_{d_{o}, k}\right)\right]  \tag{B.14}\\
& \quad=\left(d-d_{o}\right) \operatorname{tr}\left[T_{2, \sigma}\left(d_{o}, k\right)\right]+\left(d-d_{o}\right)^{2} \mathbf{o}(n / k)+\left(d-d_{o}\right) \mathbf{o}\left(n^{\epsilon+\frac{1}{2}} k^{-\frac{1}{2}}\right) .
\end{align*}
$$

Lemma B.3. For all $\beta>1$ there exists a constant $\delta>0$ such that uniformly over $\left\|\theta-\bar{\theta}_{d, k}\right\| \leq \delta_{n}$,

$$
\left|g_{n, 1}\left(\theta-\bar{\theta}_{d, k}\right)\right|=\mathbf{o}\left(n^{1 / 2-\delta} k^{-1 / 2}\right)
$$

Proof. For $u=\theta-\bar{\theta}_{d, k}$, we have

$$
g_{n, 1}(u)=-\frac{1}{2} \operatorname{tr}\left[T_{n}\left(H_{k} f_{d_{o}, \bar{\theta}_{d_{o}, k}}\right) T_{n}^{-1}\left(f_{d_{o}, \bar{\theta}_{d_{o}, k}}\right) T_{n}\left(u^{t} \nabla f_{d_{o}, \bar{\theta}_{d_{o}, k}}\right) T_{n}^{-1}\left(f_{d_{o}, \bar{\theta}_{d_{o}, k}}\right)\right] .
$$

This follows from (B.4) and Lemma below, by which $b_{\sigma}=d_{\sigma}=\frac{1}{2}$ and $c_{\sigma}=0$ (the only partition for $j=1$ being $\sigma=(\{l\})$ ). By Lemma 2.4 (supplement) $g_{n, 1}(u)$ converges to zero, but at a rate slower than $n^{1 / 2-\delta} k^{-1 / 2}$. To obtain the $\mathbf{o}\left(n^{1 / 2-\delta} k^{-1 / 2}\right)$ term, we write

$$
g_{n, 1}(u)=\Delta_{1}+\Delta_{2}+\Delta_{3},
$$

and bound the terms on the right using the other lemmas in section 2 of the supplement. We first prove that

$$
\begin{aligned}
\Delta_{1}=\operatorname{tr} & {\left[T_{n}\left(H_{k} f_{d_{o}, k}\right) T_{n}\left(f_{d_{o}, k}^{-1}\right) T_{n}\left(u^{t} \nabla f_{d_{o}, k}\right) T_{n}\left(f_{d_{o}, k}^{-1}\right)\right]-} \\
& \left(16 \pi^{4}\right) \operatorname{tr}\left[T_{n}\left(H_{k}\right) T_{n}\left(u^{t} \mathbf{c o s}\right)\right]=\mathbf{o}(1),
\end{aligned}
$$

where $\boldsymbol{\operatorname { c o s }}(x)=(1, \cos (x), \ldots, \cos (k x))$. We then prove that

$$
\Delta_{2}=\operatorname{tr}\left[T_{n}\left(H_{k}\right) T_{n}\left(u^{t} \mathbf{c o s}\right)\right]=0
$$

and finally that

$$
\begin{aligned}
\Delta_{3}= & \operatorname{tr}\left[T_{n}\left(H_{k} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right) T_{n}\left(u^{t} \nabla f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)\right] \\
& -\operatorname{tr}\left[T_{n}\left(H_{k} f_{d_{o}, k}\right) T_{n}\left(\frac{f_{d_{o}, k}^{-1}}{4 \pi^{2}}\right) T_{n}\left(u^{t} \nabla f_{d_{o}, k}\right) T_{n}\left(\frac{f_{d_{0}, k}^{-1}}{4 \pi^{2}}\right)\right] \\
= & \mathbf{o}\left(n^{1 / 2-\delta} k^{-1 / 2}\right) .
\end{aligned}
$$

To bound $\Delta_{1}$ we use Lemma 2.5 (supplement) with $b_{1}(x)=H_{k}(x), b_{2}(x)=$ $u^{t} \cos$ and $L=k^{3 / 2-\beta}$. Equation (2.6) then implies that

$$
\left|\Delta_{1}\right| \leq C \sqrt{k}\|u\| n^{\epsilon}\left(1+k^{3 / 2-\beta} k^{-1 / 2}\right)=\mathbf{o}(1)
$$

To bound Delta $_{2}$ note that for $l=0, \ldots, k$ and all $j_{1}, j_{2} \leq n$,

$$
\left(T_{n}(\cos (l x))\right)_{j_{1}, j_{2}}=\mathbb{I}_{\left|j_{1}-j_{2}\right|=l}, \quad\left(T_{n}\left(H_{k}\right)\right)_{j_{1}, j_{2}}=\sum_{j=k+1}^{n} \eta_{j} \mathbb{I}_{j=\left|j_{1}-j_{2}\right|} .
$$

Therefore

$$
\operatorname{tr}\left[T_{n}\left(H_{k}\right) T_{n}(\cos (l .)]=\sum_{j_{1}=1}^{n} \sum_{j_{2}=1}^{n} \sum_{j=k+1}^{n} \eta_{j} \mathbb{I}_{j=\left|j_{1}-j_{2}\right|} \mathbb{I}_{\left|j_{1}-j_{2}\right|=l}=0,\right.
$$

since $l \leq k$ and $j>k$. We now turn to $\Delta_{3}$. Following Lieberman et al. (2011), we consider separately the positive and negative parts of $H_{k}$ and of
$u^{t}$ cos. Hence we may treat these functions as if they were positive. We first define, for $\tilde{f}_{d_{o}, k}=\left(4 \pi^{2} f_{d_{o}, k}\right)^{-1}$,

$$
\begin{aligned}
& A_{1}=T_{n}\left(H_{k} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right), \quad B_{1}=T_{n}\left(H_{k} f_{d_{o}, k}\right) T_{n}\left(\tilde{f}_{d_{o}, k}\right), \\
& A_{2}=T_{n}\left(u^{t} \nabla f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right), \quad B_{2}=T_{n}\left(u^{t} \nabla f_{d_{o}, k}\right) T_{n}\left(\tilde{f}_{d_{o}, k}\right), \\
& \tilde{A}=T_{n}^{\frac{1}{2}}\left(H_{k} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right) T_{n}^{\frac{1}{2}}\left(u^{t} \nabla f_{d_{o}, k}\right), \\
& \tilde{B}=T_{n}^{\frac{1}{2}}\left(H_{k} f_{d_{o}, k}\right) T_{n}\left(\tilde{f}_{d_{o}, k}\right) T_{n}^{\frac{1}{2}}\left(u^{t} \nabla f_{d_{o}, k}\right), \\
& \Delta=I_{n}-T_{n}\left(f_{d_{o}, k}\right) T_{n}\left(\tilde{f}_{d_{o}, k}\right) .
\end{aligned}
$$

Using the same computations as in Lieberman et al. (2011), we find that

$$
\begin{aligned}
\left|\Delta_{3}\right| \lesssim & \left|\operatorname{tr}\left[B_{1} B_{2} \Delta\right]\right|+\left|\tilde{A}-\tilde{B} \| T_{n}^{\frac{1}{2}}\left(H_{k} f_{d_{o}, k}\right) T_{n}\left(\tilde{f}_{d_{o}, k}\right) \Delta T_{n}^{\frac{1}{2}}\left(u^{t} \nabla f_{d_{o}, k}\right)\right| \\
& +|\Delta|^{2} \sqrt{k}\|u\| n^{\epsilon} \\
\lesssim & \sqrt{k}\|u\| n^{\epsilon} k^{3 / 2-\beta}+\left|\operatorname{tr}\left[B_{1} B_{2} \Delta\right]\right| .
\end{aligned}
$$

The first term on the right is $\mathbf{o}\left(n^{1 / 2-\delta} k^{-1 / 2}\right)$. We bound the last term using Lemma 2.5 (supplement) with $b_{1}=H_{k}, b_{2}=u^{t} \boldsymbol{\operatorname { c o s }}$ and $b_{3}=1$, which implies that $\operatorname{tr}\left[B_{1} B_{2} \Delta\right]=0+O\left(\sqrt{k}\|u\| n^{\epsilon} k^{(3 / 2-\beta)_{+}}\right)=o(1)$. This achieves the proof of Lemma B.3.

Lemma B.4. Suppose that $k \leq k_{n}$ and that $l\left(f_{o}, f_{d_{o}, k}\right) \leq l_{0}^{2} \delta_{n}^{2}$. Then all elements of $\nabla_{l} l_{n}\left(d_{o}, k\right)(l=0, \ldots, k)$ are the sum of a centered quadratic form, $\mathcal{S}\left(\nabla_{l} l_{n}\left(d_{o}, k\right)\right)$ with a variance equal to $\frac{n}{2}(1+o(1))$ and a deterministic term, $\mathcal{D}\left(\nabla_{l} l_{n}\left(d_{o}, k\right)\right)$ which is $o\left(k^{(3 / 2-\beta)+} n^{\epsilon}\right)$.

Proof. For all $l=0, \ldots, k$, we have

$$
\nabla_{l} l_{n}\left(d_{o}, k\right)=\mathcal{S}\left(\nabla_{l} l_{n}\left(d_{o}, k\right)\right)+\mathcal{D}\left(\nabla_{l} l_{n}\left(d_{o}, k\right)\right),
$$

where

$$
\begin{aligned}
& \mathcal{S}\left(\nabla_{l} l_{n}\left(d_{o}, k\right)\right)= \frac{1}{2} X^{t} T_{n}^{-1}\left(f_{d_{o}, k}\right) T_{n}\left(\nabla_{l} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right) X \\
&-\frac{1}{2} \operatorname{tr}\left[T_{n}\left(f_{o}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right) T_{n}\left(\nabla_{l} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)\right], \\
& \mathcal{D}\left(\nabla_{l} l_{n}\left(d_{o}, k\right)\right)=\frac{1}{2} \operatorname{tr}\left[\left(T_{n}\left(f_{o}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)-I_{n}\right) T_{n}\left(\nabla_{l} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)\right] .
\end{aligned}
$$

Note that this is a special case of (B.2), with $j=1, b_{\sigma}=d_{\sigma}=\frac{1}{2}$ and $c_{\sigma}=0$, the only partition being $\sigma=(\{l\})$. The variance of $\mathcal{S}\left(\nabla_{l} l_{n}\left(d_{o}, k\right)\right)$ is equal
to

$$
\begin{aligned}
& \operatorname{tr}\left[\left(T_{n}\left(f_{o}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right) T_{n}\left(\nabla_{l} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)\right)^{2}\right] \\
& \quad=\frac{n}{2 \pi} \int_{-\pi}^{\pi}\left(\frac{f_{o}}{f_{d_{o}, k}}(x)\right)^{2} \cos ^{2}(l x) d x+\mathbf{O}\left(n^{\epsilon} k\right),
\end{aligned}
$$

since Lemma 2.4 (supplement) implies that the approximation error of the trace by its limiting integral is of order $O\left(n^{\epsilon}\left(k+k^{2(3 / 2-\beta) \vee 0}\right)=O\left(n^{\epsilon} k\right)\right.$. Since $\frac{f_{o}}{f_{d_{o}, k}}=e^{\Delta_{d_{o}, k}}$ (see (A.1)), the integral in the preceding equation is

$$
\begin{aligned}
& \frac{n}{2 \pi} \int_{-\pi}^{\pi}\left(1+2 \Delta_{d_{o}, k}+O\left(\Delta_{d_{o}, k}^{2}\right)\right) \cos ^{2}(l x) d x \\
& \quad=\frac{n}{2}+2 n a_{2 l}\left(d_{o}\right)+O\left(n \delta_{n}^{2}\right)=\frac{n}{2}(1+\mathbf{o}(1))
\end{aligned}
$$

where $a_{l}$ is defined at the beginning of the supplement. Lemma 1.3 (supplement) then implies that the centered quadratic form is of order $\mathbf{O}_{\mathbf{P}_{\mathbf{o}}}\left(n^{\epsilon+1 / 2}\right)$. Similarly, Lemma 2.4 (supplement) implies that

$$
\begin{aligned}
\mathcal{D}\left(\nabla_{l} l_{n}\left(d_{o}, k\right)\right) & =\frac{n}{2 \pi} \int_{-\pi}^{\pi} \frac{\left(f_{o}-f_{d_{o}, k}\right)}{f_{d_{o}, k}}(x) \cos (l x) d x+\mathbf{O}\left(n^{\epsilon}\left\|\Delta_{d_{o}, k}\right\|_{1} k\right) \\
& =\frac{n}{2 \pi} \int_{-\pi}^{\pi} \cos (l x) \Delta_{d_{o}, k}(x) d x+\mathbf{O}\left(n k^{-2 \beta}\right)+\mathbf{O}\left(n^{\epsilon} k^{3 / 2-\beta}\right) \\
& =\mathbf{O}\left(n^{\epsilon} k^{\left.(3 / 2-\beta)_{+}\right)}\right.
\end{aligned}
$$

which completes the proof of Lemma B.4.
Lemma B.5. Let $A(d)$ be the $(k+1) \times(k+1)$ matrix with entries $A_{l_{1}, l_{2}}(d)=a_{l_{1}+l_{2}}(d)$, where $a_{l}(d)=1_{l>k}\left(\theta_{o, l}-2 l^{-1}\left(d_{o}-d\right)\right)$. Suppose that $k \leq k_{n}$ and that $l\left(f_{o}, f_{d, k}\right) \leq l_{0}^{2} \delta_{n}^{2}$. Then $J_{n}(d, k)=-\left.\nabla^{2} l_{n}(d, k, \theta)\right|_{\theta=\bar{\theta}_{d, k}}$ satisfies
$\forall l_{1}, l_{2} \leq k, \quad\left|\left(J_{n}(d, k)-J_{n}\left(d_{o}, k\right)\right)_{l_{1}, l_{2}}\right|=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(\left|d-d_{o}\right| n^{\epsilon} k+S_{n}(d)\right)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(n / k)$
uniformly over $d \in\left(d_{o}-\bar{v}_{n}, d_{o}+\bar{v}_{n}\right)$ and $k \leq k_{n}$. We also have for all $l_{1}, l_{2}$

$$
\begin{equation*}
\left[J_{n}\left(d_{o}, k\right)-\frac{n}{2} I_{k+1}-\frac{n}{2} A\left(d_{o}\right)\right]_{l_{1}, l_{2}}:=n\left(R_{2 s}\right)_{l_{1}, l_{2}}+n\left(R_{2 d}\right)_{l_{1}, l_{2}}, \tag{B.16}
\end{equation*}
$$

where $\left(R_{2 s}\right)_{l_{1}, l_{2}}$ is a centered quadratic form of order $\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{-1 / 2+\epsilon}\right)$ and $\left(R_{2 d}\right)_{l_{1}, l_{2}}$ is a deterministic term of order $o\left(k n^{\epsilon-1}\right)$. For the matrix $A$, we have $\left\|A\left(d_{o}\right)\right\|=o(1)$ and $\left|A\left(d_{o}\right)\right|=O(1)$.

In particular, (B.16) implies that $\left|J_{n}\left(d_{o}, k\right)-\frac{n}{2} I_{k+1}-\frac{n}{2} A\left(d_{o}\right)\right|=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(k n^{1 / 2+\epsilon}\right)+$ $o\left(k^{2} n^{\epsilon}\right)$.

Proof. Let $d$ and $k \leq k_{n}$ be such that $l\left(f_{o}, f_{d, k}\right) \leq l_{0}^{2} \delta_{n}^{2}$ so that $d \in$ ( $d_{o}-\bar{v}_{n}, d_{o}+\bar{v}_{n}$ ) (see Corollary 3.1). Lemma B. 1 implies that for all $l_{1}, l_{2} \leq k$,

$$
\begin{aligned}
& \left(J_{n}(d, k)\right)_{l_{1}, l_{2}}-\left(J_{n}\left(d_{o}, k\right)\right)_{l_{1}, l_{2}} \\
& :=-\left(d-d_{o}\right) \sum_{\sigma \in \mathcal{S}\left(l_{1}, l_{2}\right)}\left(c_{\sigma} \operatorname{tr}\left[T_{1, \sigma}\left(d_{o}, k\right)\right]+d_{\sigma} \operatorname{tr}\left[T_{2, \sigma}\left(d_{o}, k\right)\right]\right)+\mathbf{O}_{\mathbf{P}_{\mathbf{o}}}\left(S_{n}(d)\right) .
\end{aligned}
$$

Lemma 2.4 (supplement) implies that

$$
\operatorname{tr}\left[T_{i, \sigma}\left(d_{o}, k\right)\right]=O\left(k n^{\epsilon}\right), \quad i=1,2
$$

so that (B.15) is satisfied since this term is $\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{1-\delta} / k\right)$. We then use expression (B.2), with $\sigma \in\{(\{1\},\{2\}),(\{1,2\})\}$ and we denote $\sigma_{1}$ and $\sigma_{2}$ the first and the second partition respectively. Note that $c_{\sigma_{1}}=d_{\sigma_{2}}=1 / 2$, $c_{\sigma_{2}}=0$ and $d_{\sigma_{1}}=1$. From Lemma 2.4 (supplement), the quadratic form in $\left(J_{n}\left(d_{o}, k\right)\right)_{l_{1}, l_{2}}$ is associated to a matrix whose Frobenius-norm is $O(\sqrt{n})$ and whose spectral norm is $O\left(n^{\epsilon}\right)$. Hence, this quadratic form is $\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{1 / 2+\epsilon}\right)$. Also by Lemma 2.4 (supplement), the deterministic terms can be written as

$$
\begin{aligned}
& \frac{n}{4 \pi} \int_{-\pi}^{\pi} \cos \left(l_{1} x\right) \cos \left(l_{2} x\right)\left(1+\Delta_{d_{o}, k}\right)(x) d x+o\left(k n^{\epsilon}\right) \\
& \quad=\frac{n}{2}\left(\mathbb{I}_{l_{1}=l_{2}}+a_{l_{1}+l_{2}}\left(d_{o}\right)\right)+o\left(k n^{\epsilon}\right),
\end{aligned}
$$

and Lemma B. 5 is proved.

## APPENDIX C: PROOF OF LEMMA 3.5

Under the conditions of Theorem 2.1 we have $k=k_{n}$ and $\beta>1$, and we may assume (by Lemma 3.2) that $l\left(f_{o}, f_{d, k, \theta}\right) \leq l_{0}^{2} \delta_{n}^{2}$. Fixing $d$ and $k$, we develop $\theta \rightarrow l_{n}(d, k, \theta)$ in $\bar{\theta}_{d, k}$. From Lemma B. 1 in Appendix B it follows that

$$
\begin{align*}
l_{n}(d, k, \theta)-l_{n}(d, k)= & \sum_{j=1}^{J} \frac{\left(\theta-\bar{\theta}_{d, k}\right)^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!} \\
& +\left(d-d_{o}\right) \sum_{j=2}^{J} \frac{g_{n, j}\left(\theta-\bar{\theta}_{d, k}\right)}{j!}+S_{n}(d), \tag{C.1}
\end{align*}
$$

where $S_{n}(d)$ is as in (B.5). Substituting (C.1) in the definition of $I_{n}(d, k)$ in (3.1), we obtain

$$
\begin{align*}
& I_{n}(d, k)=\int_{\left\|\theta-\bar{\theta}_{d, k}\right\| \leq 2 l_{0} \delta_{n}} e^{l_{n}(d, k, \theta)-l_{n}(d, k)} d \pi_{\theta \mid k}(\theta)  \tag{C.2}\\
& =e^{S_{n}(d)} \pi_{\theta \mid k}\left(\bar{\theta}_{d, k}\right) \int_{\|u\| \leq 2 l_{0} \delta_{n}} e^{\sum_{j=1}^{J} \frac{1}{j!} u^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)+\left(d-d_{o}\right) \sum_{j=2}^{J} \frac{g_{n, j}(u)}{j!}+h_{k} u} d u \\
& =e^{S_{n}(d)} \pi_{\theta \mid k}\left(\bar{\theta}_{d_{o}, k}\right) \int_{\|u\| \leq 2 l_{0} \delta_{n}} e^{\sum_{j=1}^{J} \frac{1}{j!} u^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)+\left(d-d_{o}\right) \sum_{j=2}^{J} \frac{g_{n, j}(u)}{j!}+h_{k} u} d u .
\end{align*}
$$

The first equality follows from the definition of $I_{n}(d, k)$ and Lemma 3.3, by which we may replace the domain of integration by $\left\{\theta:\left\|\theta-\bar{\theta}_{d, k}\right\| \leq 2 l_{0} \delta_{n}\right\}$. The second equality follows from the assumptions on $\pi_{\theta \mid k}$ in prior $\mathbf{A}$, the transformation $u=\theta-\bar{\theta}_{d, k}$ and substitution of (C.1). Also the third equality follows from the assumptions on $\pi_{\theta \mid k}$ : these imply that

$$
\left|\log \pi_{\theta \mid k}\left(\bar{\theta}_{d, k}\right)-\log \pi_{\theta \mid k}\left(\bar{\theta}_{d_{o}, k}\right)\right|=\left|d_{o}-d\right|\left|h_{k}^{t} \eta_{[k]}\right|+o(1)=O\left(\left|d-d_{o}\right|(n / k)^{\frac{1}{2}-\epsilon}\right)+o(1)
$$

for some $\epsilon>0$. Thus, the factor $e^{S_{n}(d)} \pi_{\theta \mid k}\left(\bar{\theta}_{d, k}\right)$ on the second line of (C.2) may be replaced by $e^{S_{n}(d)} \pi_{\theta \mid k}\left(\bar{\theta}_{d_{o}, k}\right)$. Because $S_{n}\left(d_{o}\right)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)$, (C.2) implies that

$$
\begin{equation*}
I_{n}\left(d_{o}, k\right)=\left(1+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)\right) \int_{\|u\| \leq 2 l_{0} \delta_{n}} \exp \left\{h_{k} u+\sum_{j=1}^{J} \frac{u^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}\right\} d u \tag{C.3}
\end{equation*}
$$

The most involved part of the proof is to establish the bounds

$$
\begin{align*}
& \pi_{\theta \mid k}\left(\bar{\theta}_{d_{o}, k}\right) \int_{\|u\| \leq l_{0} \delta_{n}} \exp \left\{h_{k} u+\sum_{j=1}^{J} \frac{u^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}\right\} d u  \tag{C.4}\\
& \leq I_{n}(d, k)=e^{S_{n}(d)} \pi_{\theta \mid k}\left(\bar{\theta}_{d_{o}, k}\right) \int_{\|u\| \leq 2 l_{0} \delta_{n}} e^{\sum_{j=1}^{J} \frac{1}{j!} u^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)+\left(d-d_{o}\right) \sum_{j=2}^{J} \frac{g_{n, j}(u)}{j!}+h_{k} u} d u \\
& \leq \pi_{\theta \mid k}\left(\bar{\theta}_{d_{o}, k}\right) \int_{\|u\| \leq 3 l_{0} \delta_{n}} \exp \left\{h_{k} u+\sum_{j=1}^{J} \frac{u^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}\right\} d u .
\end{align*}
$$

Since the posterior distribution of $\theta$ conditional on $k=k_{n}$ and $d=d_{o}$ concentrates at $\bar{\theta}_{d_{o}, k}$ at a rate bounded by $l_{0} \delta_{n}$ (this follows from Lemma
3.2, with the restriction to $d=d_{o}$ ), the left- and right-hand side of (C.4) are asymptotically equal, up to a factor ( $1+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)$ ). By (C.3), the left- and right-hand side are actually equal to $I_{n}\left(d_{o}, k\right)$. This implies that $I_{n}(d, k)=$ $e^{S_{n}(d)} I_{n}\left(d_{o}, k\right)$, which is the required result.

In the remainder we prove (C.4). To do so we construct below a change of variables $v=\psi(u)$, which satisfies

$$
\begin{align*}
h_{k} v+\sum_{j=1}^{J} \frac{v^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}= & h_{k} u+\sum_{j=1}^{J} \frac{u^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!} \\
& +\left(d-d_{o}\right) \sum_{j=2}^{J} \frac{g_{n, j}(u)}{j!}+O\left(S_{n}(d)\right) \tag{C.5}
\end{align*}
$$

for all $\|u\| \leq 2 l_{0} \delta_{n}$. We first define the notation required in the definition of $\psi$ in (C.8) below. Recall from (B.4) in Lemma B. 1 that $g_{n, j}(u)$ can be decomposed as

$$
g_{n, j}(u)=n \sum_{\sigma \in \mathcal{S}_{j}}\left(c_{\sigma}-d_{\sigma}\right) \sum_{l_{1}, \ldots, l_{j}=0}^{k} u_{l_{1}} \ldots u_{l_{j}} g_{l_{1}, \ldots, l_{j}}^{(j)}, \quad j=2, \ldots J,
$$

where $g_{l_{1}, \ldots, l_{j}}^{(j)}$ depends on $\sigma$. For ease of presentation however we omit this dependence in the notation. Using Lemma 2.4 (supplement) and (B.4) in Lemma B.1, it follows that for all $j \geq 2$ and $\left(l_{1}, \ldots, l_{j}\right) \in\{0, \ldots, k\}^{j}$,

$$
\begin{gather*}
g_{l_{1}, \ldots, l_{j}}^{(j)}=\gamma_{l_{1}, \ldots, l_{j}}^{(j)}+r_{l_{1}, \ldots, l_{j}}^{(j)},  \tag{C.6}\\
\gamma_{l_{1}, \ldots, l_{j}}^{(j)}=\frac{1}{2 \pi} \int_{-\pi}^{\pi} H_{k}(x) \cos \left(l_{1} x\right) \cdots \cos \left(l_{j} x\right) d x
\end{gather*}
$$

Let $\bar{G}^{(2)}$ denote the matrix with elements $\gamma_{l_{1}, l_{2}}^{(2)}$, and $G^{(2)}$ the matrix with elements $g_{l_{1}, l_{2}}^{(2)}$. By direct calculation it follows that

$$
\begin{equation*}
\gamma_{l_{1}, l_{2}}^{(2)}=\mathbb{I}_{l_{1}+l_{2}>k} \frac{1}{2\left(l_{1}+l_{2}\right)} . \tag{C.7}
\end{equation*}
$$

Similarly, for all $j \geq 3$ and $l_{1}, \ldots, l_{j} \in\{0,1, \ldots, k\}$ we define

$$
\left(G^{(j)}(u)\right)_{l_{1}, l_{2}}=\sum_{l_{3}, \ldots, l_{j}=0}^{k} g_{l_{1}, \ldots, l_{j}}^{(j)} u_{l_{3}} \ldots u_{l_{j}},\left(\bar{G}^{(j)}(u)\right)_{l_{1}, l_{2}}=\sum_{l_{3}, \ldots, l_{j}=0}^{k} \gamma_{l_{1}, \ldots, l_{j}}^{(j)} u_{l_{3}} \ldots u_{l_{j}} .
$$

In contrast to $G^{(2)}$ and $\bar{G}^{(2)}, G^{(j)}(u)$ and $\bar{G}^{(j)}(u)$ depend on $u$. For notational convenience we will also write $G^{(2)}(u)$ and $\bar{G}^{(2)}(u)$. Finally, let $\tilde{I}_{k}=J_{n}\left(d_{o}, k\right) / n$ be the normalized Fisher information.

We now define the transformation $\psi$ :

$$
\begin{gather*}
\psi(u)=\left(I_{k+1}-\left(d-d_{o}\right) D(u)\right) u \text {, with }  \tag{C.8}\\
D(u)=\left(\tilde{I}_{k}+L(u)\right)^{-1} G^{t}(u), \quad G(u)=\sum_{j=2}^{J} \frac{1}{j!} \sum_{\sigma \in \mathcal{S}_{j}}\left(c_{\sigma}-d_{\sigma}\right) G^{(j)}(u), \\
9) \quad(L(u))_{l_{1}, l_{2}}=-\sum_{j=3}^{J} \frac{1}{n(j-1)!} \sum_{l_{3}, \ldots, l_{j}=0}^{k} u_{l_{3}} \ldots u_{l_{j}} \nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right) .
\end{gather*}
$$

The construction of $G(u)$ is such that

$$
\begin{equation*}
n u^{t} G(u) u=\sum_{j=2}^{J} g_{n, j}(u) . \tag{C.10}
\end{equation*}
$$

Analogous to $G(u)$ and $D(u)$ we define $\bar{G}(u)=\sum_{j=2}^{J} \frac{1}{j!} \sum_{\sigma \in \mathcal{S}_{j}}\left(c_{\sigma}-d_{\sigma}\right) \bar{G}^{(j)}(u)$ and $\bar{D}(u)=\left(\tilde{I}_{k}+L(u)\right)^{-1} \bar{G}^{t}(u)$. After substitution of $v=\psi(u)$, and using (C.25) in Lemma C. 1 it follows that

$$
\begin{aligned}
& \sum_{j=3}^{J} \frac{\left(v^{(j)}-u^{(j)}\right) \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}= \\
& \quad-\left(d-d_{o}\right) \sum_{j=3}^{J} \sum_{l_{1}, \ldots, l_{j}=0}^{k}(D(u) u)_{l_{1}} u_{l_{2}} \ldots u_{l_{j}} \frac{\nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right)}{(j-1)!}+\mathbf{O}\left(S_{n}(d)\right) .
\end{aligned}
$$

The definitions of $D(u)$ and $L(u)$ and (C.10) imply that

$$
\begin{aligned}
& -n(v-u)^{t} \tilde{I}_{k} u=n\left(d-d_{o}\right) u^{t} D^{t}(u) \tilde{I}_{k} u=n\left(d-d_{o}\right) u^{t} G(u)\left(\tilde{I}_{k}+L(u)\right)^{-1} \tilde{I}_{k} u \\
& =n\left(d-d_{o}\right) u^{t} G(u)\left(I_{k+1}-\left(\tilde{I}_{k}+L(u)\right)^{-1} L(u)\right) u \\
& =\left(d-d_{o}\right) \sum_{j=2}^{J} \frac{1}{j!} g_{n, j}(u)-n\left(d-d_{o}\right)(D(u) u)^{t} L(u) u \\
& =\left(d-d_{o}\right) \sum_{j=2}^{J} \frac{1}{j!} g_{n, j}(u)-\left(d-d_{o}\right) \sum_{j=3}^{J-1} \sum_{l_{1}, \ldots, l_{j}=0}^{k}(D(u) u)_{l_{1}} u_{l_{2}} \ldots u_{l_{j}} \frac{\nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right)}{(j-1)!} .
\end{aligned}
$$

At the same time, the definition of $\tilde{I}_{k}$ implies that

$$
\frac{1}{2}\left(v^{(2)}-u^{(2)}\right) \nabla^{2} l_{n}\left(d_{o}, k\right)=-n(v-u)^{t} \tilde{I}_{k} u-\frac{n(v-u)^{t} \tilde{I}_{k}(v-u)}{2}
$$

Combining the preceding results, we find that

$$
\begin{aligned}
& h_{k} v+\sum_{j=1}^{J} \frac{v^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}-\left(h_{k} u+\sum_{j=1}^{J} \frac{u^{(j)} \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}+\left(d-d_{o}\right) \sum_{j=2}^{J} \frac{g_{n, j}(u)}{j!}\right) \\
& \quad=h_{k}(v-u)+(v-u)^{t} \nabla l_{n}\left(d_{o}, k\right)-\frac{n(v-u)^{t} \tilde{I}_{k}(v-u)}{2}+O\left(S_{n}(d)\right) \\
& \quad=(v-u)^{t} \nabla l_{n}\left(d_{o}, k\right)+O\left(S_{n}(d)\right),
\end{aligned}
$$

where the last equality follows from (C.24) below in Lemma C.1, together with the assumption on $h_{k}$ in prior $\mathbf{A}$ in (2.2).

Apart from the term $(v-u)^{t} \nabla l_{n}\left(d_{o}, k\right)$ on the last line, the preceding display implies (C.5). Hence, to complete the proof of (C.5) it suffices to show that

$$
\begin{equation*}
(v-u)^{t} \nabla l_{n}\left(d_{o}, k\right)=-\left(d-d_{o}\right) u^{t} D^{t}(u) \nabla l_{n}\left(d_{o}, k\right)=O\left(S_{n}(d)\right) . \tag{C.11}
\end{equation*}
$$

The proof of (C.11) consists of the following steps:

$$
\begin{align*}
\left|u^{t}(D(u)-\bar{D}(u))^{t} \nabla l_{n}\left(d_{o}, k\right)\right| & =\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{\frac{1}{2}-\delta} k^{-\frac{1}{2}}\right),  \tag{C.12}\\
\left(d-d_{o}\right) u^{t} \bar{D}^{t}(u) \mathcal{D}\left(\nabla l_{n}\left(d_{o}, k\right)\right) & =O\left(S_{n}(d)\right),  \tag{C.13}\\
\left(d-d_{o}\right) u^{t} \bar{D}^{t}(u) \mathcal{S}\left(\nabla l_{n}\left(d_{o}, k\right)\right) & =O\left(S_{n}(d)\right), \tag{C.14}
\end{align*}
$$

where $\mathcal{S}\left(\nabla l_{n}\left(d_{o}, k\right)\right)$ denotes the centered quadratic form in $\nabla l_{n}\left(d_{o}, k\right)$, and $\mathcal{D}\left(\nabla l_{n}\left(d_{o}, k\right)\right)$ the remaining deterministic term. We will use the same notation below for $L(u)$.

Equation (C.12) follows from Lemma B. 4 and (C.22) in Lemma C. 1 below, which imply that the left-hand side equals $\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left((\sqrt{k}\|u\|)^{2} n^{-1+\epsilon} k \sqrt{n}\right)=$ $\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{\frac{1}{2}-\delta} k^{-\frac{1}{2}}\right)$, for some $\delta>0$. For the proof of (C.13), note that Lemma B. 4 implies

$$
\left|u^{t} \bar{D}^{t}(u) \mathcal{D}\left(\nabla l_{n}\left(d_{o}, k\right)\right)\right| \lesssim\|\bar{D}(u) u\| \sqrt{k} k^{(3 / 2-\beta)_{+}+\epsilon} .
$$

Combined with Lemma C. 1 , this implies that the left-hand side is $O\left(\sqrt{k} k^{5 / 2-2 \beta+\epsilon}\right)$, which is $O\left(S_{n}(d)\right)$. The proof of (C.14) is more involved. Recall that $\bar{D}(u)$ is defined as $\bar{D}(u)=\left(\tilde{I}_{k}+L(u)\right)^{-1} \bar{G}^{t}(u)$. Using (B.16) in Lemma B.5, we obtain

$$
\left(\tilde{I}_{k}+L(u)\right)^{-1}=2\left[I_{k+1}-\left(A\left(d_{o}\right)+R_{2 s}+R_{2 d}+L(u)\right)\left(1+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)\right)\right] .
$$

Substituting this in $\bar{D}(u)$, it follows that (C.14) can be proved by controlling $\bar{G}^{(j)} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right), \bar{G}^{(j)} A\left(d_{o}\right) \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right), \bar{G}^{(j)} R_{2 d} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right), \bar{G}^{(j)} \mathcal{D}(L(u)) \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)$, $\bar{G}^{(j)} R_{2 s} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)$ and $\bar{G}^{(j)} \mathcal{S}(L(u)) \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)$ for all $j=3, \ldots, J$. To do so, first note that Lemma B. 5 implies that $\left\|\bar{G}^{(j)} R_{2 s} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)\right\|=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{\epsilon} \sqrt{k}\right)$. Hence,

$$
\left|u^{t} \bar{G}^{(j)} R_{2 s} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)\right|=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(k^{-\beta+1+\epsilon}\right)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)
$$

which clearly is $O\left(S_{n}(d)\right)$. The terms $\bar{G}^{(j)} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right), \bar{G}^{(j)} A\left(d_{o}\right) \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)$, $\bar{G}^{(j)} R_{2 d} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)$ and $\bar{G}^{(j)} \mathcal{D}(L(u)) \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)$ can be written as quadratic forms $Z^{t} M Z-\operatorname{tr}[M]$, where, for a sequence $\left(b_{l}\right)_{l=0}^{k}$ and a function $g$ with $\|g\|_{\infty}<\infty, M$ is of the form

$$
T_{n}^{\frac{1}{2}}\left(f_{o}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right) T_{n}\left(g(x) f_{d_{o}, k}(x) \sum_{l} b_{l} \cos (l x)\right) T_{n}^{-1}\left(f_{d_{o}, k}\right) T_{n}^{\frac{1}{2}}\left(f_{o}\right)
$$

$Z$ being a vector of $n$ independent standard Gaussian random variables. Using Lemma 2.4 (supplement) it can be seen that $|M|^{2} \leq n\left(\sum_{l} b_{l}^{2}+k / n\right)$. Lemma 1.3 (supplement) with $\alpha=\epsilon+1 / 2$ then implies that

$$
\begin{equation*}
P_{o}\left(\left|Z^{t} M Z-\operatorname{tr}[M]\right|>n^{\epsilon+\frac{1}{2}}\left(\sum_{l} b_{l}^{2}+\frac{k}{n}\right)^{\frac{1}{2}}\right) \leq e^{-c n^{\epsilon}} \tag{C.15}
\end{equation*}
$$

For all $j \in\{3, \ldots, J\}$, the four terms above can now be bounded for a particular choice of $g$ and $b_{l}$.

- Bound on $\bar{G}^{(j)} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)$. For all $l_{2}, \ldots, l_{j} \in\{0, \ldots, k\}$, set $b_{l}=$ $\gamma_{l_{2}, l, l_{4}, \ldots, l_{j}}^{(j)}$ and $g(x)=1$. Then we have

$$
\left.\sum_{l=0}^{k} b_{l} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)\right)_{l_{1}}=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{\frac{1}{2}+\epsilon}\left(\sum_{l} b_{l}^{2}\right)^{\frac{1}{2}}\right)
$$

By induction it can be shown that

$$
\begin{align*}
& \frac{1}{2 \pi} \int_{-\pi}^{\pi} \cos \left(l_{0} x\right) \cos \left(l_{1} x\right) \cos \left(l_{2} x\right) \cdots \cos \left(l_{j} x\right) d x \\
& \quad=2^{-j} \sum_{\epsilon_{1}, \ldots, \epsilon_{j} \in\{-1,1\}} \mathbb{I}_{l_{0}+\sum_{i=1}^{j} \epsilon_{i} l_{i}=0} \tag{C.16}
\end{align*}
$$

Consequently, $\sum_{l} b_{l}^{2}=O\left(k^{-1}\right)$ for all $l_{2}, l_{4}, \ldots, l_{j} \in\{0, \ldots, k\}$. Using the fact that $\sum_{l}\left|u_{l}\right|=o(1)$, we obtain that $\left(\bar{G}^{(j)} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)\right)_{l_{2}}=$ $\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{1 / 2+\epsilon} k^{-1 / 2}\right)$, for all $l_{2} \in\{0, \ldots, k\}$. This implies that

$$
\begin{equation*}
\left\|\bar{G}^{(j)} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)\right\|=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{1 / 2+\epsilon}\right)=O\left(S_{n}(d)\right) . \tag{C.17}
\end{equation*}
$$

- Bound on $A\left(d_{o}\right) \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)$. Set $b_{l}=\mathbb{I}_{l+l_{1} \geq k} \theta_{o, l+l_{1}}$ and $g(x)=1$, then

$$
\left(A\left(d_{o}\right) \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)\right)_{l_{1}}=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{1 / 2+\epsilon} k^{-\beta}\right), \quad \forall l_{1} \in\{0, \ldots, k\}
$$

Combined with Lemma C. 1 this implies that

$$
\left\|\bar{G}^{(j)} A\left(d_{o}\right) \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)\right\|=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{1 / 2+\epsilon} k^{-\beta+1 / 2}\right)
$$

- Bound on $\bar{G}^{(j)} R_{2 d} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right.$. Set $b_{l}=\left(R_{2 d}\right)_{l_{1}, l}$ and $g(x)=1$, for all $l_{1}=0, \ldots, k$, then Lemmas B. 5 and C. 1 lead to

$$
\left\|\bar{G}^{(j)} R_{2 d} \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)\right\|=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{1 / 2+\epsilon} k^{2} n^{-1}\right)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{-1 / 2+\epsilon} k^{2}\right)
$$

- Bound on $\bar{G}^{(j)} \mathcal{D}(L(u)) \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)$. For all $l_{1}, l_{3}, \ldots, l_{j}, l_{3}^{\prime}, \ldots, l_{j^{\prime}}^{\prime} \in\{0, \ldots, k\}$, set

$$
\begin{aligned}
b_{l}= & \frac{1}{n} \operatorname{tr}\left[T_{n}^{-1}\left(f_{d_{o}}\right) T_{n}\left(\sum_{l_{2}=0}^{k} \gamma_{l_{1}, l_{2}, \ldots, l_{j}}^{(j)} \cos \left(l_{2} x\right) g_{1}(x) f_{d_{o}}(x)\right) \times\right. \\
& \left.T_{n}^{-1}\left(f_{d_{o}}\right) T_{n}\left(\cos (l x) g_{2}(x) f_{d_{o}}(x)\right) \cdots T_{n}^{-1}\left(f_{d_{o}}\right) T_{n}\left(g_{r}(x) f_{d_{o}}(x)\right)\right]
\end{aligned}
$$

where $g_{1}(x), \ldots, g_{r}(x)$ are products of functions of the form $\cos \left(l_{i}^{\prime} x\right)$ and $g_{1}(x) \ldots g_{r}(x)=\cos \left(l_{3}^{\prime} x\right) \ldots \cos \left(l_{j^{\prime}}^{\prime} x\right)$. Lemmas 2.1 and 2.6 in the supplement, together with (C.16), imply that

$$
\sum_{l} b_{l}^{2}=O\left(k^{-1}\right), \quad \text { and } \quad\left\|\bar{G}^{(j)} \mathcal{D}(L(u)) \mathcal{S}\left(\nabla l_{n}\left(d_{o}\right)\right)\right\|=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{\epsilon+1 / 2}\right)
$$

Consequently, the contribution to all these terms in $(v-u)^{t} \nabla l_{n}\left(d_{o}, k\right)$ is of order $\mathbf{O}\left(S_{n}(d)\right)$.

We control $u^{t} \bar{G}^{(j)} \mathcal{S}(L(u)) \mathcal{S}\left(\nabla l_{n}\left(d_{o}, k\right)\right)$, by bounding $\left\|\bar{G}^{(j)} \mathcal{S}(L(u))\right\|$ using a similar idea. Indeed, for all $l_{1}, l_{2} \leq k,\left(\bar{G}^{(j)} \mathcal{S}(L(u))\right)_{l_{1}, l_{2}}$ can be written as a sum of terms of the form $\left(Z^{t} M_{l_{1}, l_{2}} Z-\operatorname{tr}\left(M_{l_{1}, l_{2}}\right)\right) / n$, where $Z$ is a vector of $n$ independent standard Gaussian random variables, and $M_{l_{1}, l_{2}}$ has the form

$$
\begin{aligned}
& T_{n}^{\frac{1}{2}}\left(f_{o}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)\left(\prod_{i<i_{0}} T_{n}\left(\nabla_{\sigma}(i) f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)\right) \times \\
& T_{n}\left(\sum_{l=0}^{k} \gamma_{l, l_{1}, l_{2}, \ldots, l_{j-1}}^{(j)} \cos (l x) \nabla_{\sigma\left(i_{0}\right)-\{l\}} f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right) \prod_{i<i_{0}} T_{n}\left(\nabla_{\sigma}(i) f_{d_{o}, k}\right) T_{n}^{\frac{1}{2}}\left(f_{o}\right)
\end{aligned}
$$

We can use the same argument as in (C.15) since for all $l_{1}, l_{2}, \ldots, l_{j-1}$

$$
\begin{aligned}
\left|M_{l_{1}, l_{2}}\right| & \lesssim \left\lvert\, T_{n}^{-\frac{1}{2}}\left(\left.f_{d_{o}, k} T_{n}\left(\sum_{l=0}^{k} \gamma_{l, l_{1}, l_{2}, \ldots, l_{j-1}}^{(j)} \cos (l x) \nabla_{\sigma\left(i_{0}\right)-\{l\}} f_{d_{o}, k}\right) T_{n}^{-\frac{1}{2}}\left(f_{d_{o}, k}\right) \right\rvert\,\right.\right. \\
& =O\left(n^{1 / 2+\epsilon} k^{-1 / 2}\right)
\end{aligned}
$$

Hence, it follows that $n^{-1}\left[Z^{t} M_{l_{1}, l_{2}} Z-\operatorname{tr}\left(M_{l_{1}, l_{2}}\right)\right]=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{-1 / 2+\epsilon} k^{-1 / 2}\right)$ and

$$
\begin{equation*}
u^{t} \bar{G} \mathcal{S}(L(u)) \mathcal{S}\left(\nabla l_{n}\left(d_{o}, k\right)\right)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(\|u\| n^{\epsilon} k\right)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{1 / 2-\delta} k^{-1 / 2}\right) \tag{C.19}
\end{equation*}
$$

Combining (C.19) and (C.17)-(C.18), we obtain (C.14). This in turn finishes the proof of (C.11), since

$$
(v-u)^{t} \nabla l_{n}\left(d_{o}, k\right)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(\left|d-d_{o}\right| n^{\frac{1}{2}-\delta} k^{-\frac{1}{2}}\right)=O\left(S_{n}(d)\right) .
$$

We now prove that $\psi(u)$ is a one-to-one transformation. First note that $\psi(u)$ is continuously differentiable for all $\|u\| \leq 2 l_{0} \delta_{n}$. This follows from the definition $\psi(u)=\left(I_{k+1}-\left(d-d_{o}\right)\left(\tilde{I}_{k}+L(u)\right)^{-1} G^{t}(u)\right) u$, the fact that $G(u)$ and $L(u)$ are polynomial in $u$ and Lemma C.1, by which $\|L(u)\|=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)$. To prove that $\psi(u)$ is also one-to-one, we bound the spectral norm of the Jacobian

$$
\psi^{\prime}(u)=I_{k+1}-\left(d-d_{o}\right) D(u)-\left(d-d_{o}\right)\left(D^{\prime}(u) u\right)
$$

where $\left(D^{\prime}(u) u\right)$ is the $(k+1) \times(k+1)$ matrix with elements

$$
\sum_{l=0}^{k} u_{l} \frac{\partial(D(u))_{l_{1}, l}}{\partial u_{l_{2}}}, \quad l_{1}, l_{2}=0, \ldots, k
$$

For $\psi(u)$ to be one-to-one, it suffices to have $\psi^{\prime}(u)=I_{k+1}\left(1+\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1)\right)$.
By (C.24) in Lemma C. 1 below, we have $\left|d-d_{o}\right|\|D(u)\|=\mathbf{O}_{P_{o}}\left(\left|d-d_{o}\right|\right)$. Therefore we only need to control the spectral norm of $D^{\prime}(u) u$. For all $l_{1}, l_{2}$, we have
$\left(D^{\prime}(u) u\right)_{l_{1}, l_{2}}=\left[-\left(\tilde{I}_{k}+L(u)\right)^{-1} \frac{\partial L(u)}{\partial u_{l_{2}}}\left(\tilde{I}_{k}+L(u)\right)^{-1} G^{t}(u) u+\left(\tilde{I}_{k}+L(u)\right)^{-1} \frac{\partial G^{t}(u)}{\partial u_{l_{2}}} u\right]_{l_{1}}$.
Both $(G(u))_{l_{1}, l_{2}}$ and $(L(u))_{l_{1}, l_{2}}$ can be written as

$$
F_{l_{1}, l_{2}}(u ; \tau, b):=\sum_{j=2}^{J} \tau_{j} \sum_{l_{3}, \ldots, l_{j}=0}^{k} u_{l_{3}} \cdots u_{l_{j}} b_{l_{1}, l_{2}, \ldots, l_{j}}
$$

where the constants $\tau_{j}, b_{l_{1}, \ldots, l_{j}}$ are different for $G$ and $L$, and $b$ is symmetric in its indices. In particular, $\tau_{2}=0$ in the case of $L$. Using this generic notation for $G(u)$ and $L(u)$, we find that for all $v \in \mathbb{R}^{k+1}$ and all $l_{1}, l_{2} \leq k$,

$$
\left(\frac{\partial F(u ; \tau, b) v}{\partial u_{l_{2}}}\right)_{l_{1}}=\sum_{j=3}^{J} \tau_{j}(j-3+1) \sum_{l_{3}, \ldots, l_{j}=0}^{k} v_{l_{3}} u_{l_{4}} \cdots u_{l_{j}} b_{l_{1}, l_{2}, \ldots, l_{j}}:=F\left(v, u ; \tau^{\prime}, b\right)
$$

where $\tau_{j}^{\prime}=\tau_{j}(j-3+1), j=3, \ldots, J$. It therefore has the same form as $F\left(u ; \tau^{\prime}, b\right)$, with $v$ replacing one of the $u$ 's. Applying this to the first term of (C.20), with $v=\left(\tilde{I}_{k}+L(u)\right)^{-1} G^{t}(u) u$, we find that

$$
\left|\left(\tilde{I}_{k}+L(u)\right)^{-1} F\left(v, u ; \tau^{\prime}, b\right)\right| \lesssim\left|F\left(v, u ; \tau^{\prime}, b\right)\right|=\mathbf{O}(1)
$$

where we used (C.21) and (C.24) from Lemma C.1. The second term of (C.20) is treated similarly with $v=u$ so that we finally obtain

$$
\left|D^{\prime}(u) u\right|=\mathbf{O}(1),
$$

and $\psi$ is one-to-one on $\left\{u:\|u\| \leq 2 l_{0} \delta_{n}\right\}$. Using the above bounds we also deduce that the Jacobian is equal to $\exp \left(O\left(S_{n}(d)\right)\right)$, since

$$
\begin{aligned}
& \log \operatorname{det}[\mathrm{Jac}]=\log \operatorname{det}\left[I_{k+1}-\left(d-d_{o}\right) D(u)-\left(d-d_{o}\right) D^{\prime}(u) u\right] \\
& \left.=O\left[\left(d-d_{o}\right)\left(|\operatorname{tr}[D(u)]|+\operatorname{tr}\left[D^{\prime}(u) u\right]\right)+\left(d-d_{o}\right)^{2}\left(|D(u)|^{2}+\left|D^{\prime}(u) u\right|^{2}\right)\right)\right] \\
& =O\left(\sqrt{k}\left(d-d_{o}\right)+k\left(d-d_{o}\right)^{2}\right)=O\left(S_{n}(d)\right)
\end{aligned}
$$

This finishes the proof of (C.4), and hence the proof of Lemma 3.5.
Lemma C.1. Let $v=\psi(u)$, with $\psi$ as in (C.8). Under the conditions of Lemma 3.5, we have

$$
\begin{align*}
|L(u)| & =\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{-1 / 2+\epsilon} k\right)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1),  \tag{C.21}\\
|G-\bar{G}| & =\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{-1 / 2+\epsilon} k\right)=\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}(1),  \tag{C.22}\\
|D(u)| & =\mathbf{O}_{P_{o}}(1),  \tag{C.23}\\
\|u-\psi(u)\| & \lesssim\left|d-d_{o}\right| \mathbf{O}_{P_{o}}(\|u\|), \tag{C.24}
\end{align*}
$$

and

$$
\begin{align*}
& \sum_{j=3}^{J} \frac{\left(v^{(j)}-u^{(j)}\right) \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}  \tag{C.25}\\
& =-\left(d-d_{o}\right) \sum_{j=3}^{J} \sum_{l_{1}, \ldots, l_{j}=0}^{k}(D(u) u)_{l_{1}} u_{l_{2}} \ldots u_{l_{j}} \frac{\nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right)}{(j-1)!}+O\left(S_{n}(d)\right),
\end{align*}
$$

uniformly over $\|u\| \leq 2 l_{0} \delta_{n}$.

Proof. We first prove (C.21). From (B.2), we recall that $\nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right)$ is the sum of a centered quadratic form $\mathcal{S}\left(\nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right)\right)$ and a deterministic term $\mathcal{D}\left(\nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right)\right)$. For all $l_{1}, \ldots, l_{j}, \mathcal{S}\left(\nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right)\right)$ equals

$$
X^{t}\left(T_{n}^{-1}\left(f_{d_{o}, k}\right) \sum_{\sigma \in \mathcal{S}_{j}} b_{\sigma} B_{\sigma}\left(d_{o}\right)\right) X-\operatorname{tr}\left[T_{n}\left(f_{o}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right) \sum_{\sigma \in \mathcal{S}_{j}} b_{\sigma} B_{\sigma}\left(d_{o}\right)\right],
$$

with $B_{\sigma}\left(d_{o}\right):=B_{\sigma}\left(d_{o}, \bar{\theta}_{d_{o}, k}\right)$ as defined in (B.1). Using Lemma 1.3 (supplement) together with (B.10) we obtain that for all $l_{1}, \ldots, l_{j}, \mathcal{S}\left(\nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right)\right)=$ $\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(n^{\frac{1}{2}+\epsilon}\right)$, and its contribution to $|L(u)|$ is $\mathbf{o}_{\mathbf{P}_{\mathbf{o}}}\left(k(\sqrt{k}\|u\|)^{j-2} n^{-\frac{1}{2}+\epsilon}\right)=$ $\mathbf{o}_{\mathbf{P}}\left(n^{-1 / 2-\delta} k\right)$. The deterministic term in (B.2) is
$\mathcal{D}\left(\nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right)\right)=\sum_{\sigma} c_{\sigma} \operatorname{tr}\left[B_{\sigma}\left(d_{o}\right)\right]+\sum_{\sigma} d_{\sigma} \operatorname{tr}\left[\left(T_{n}\left(f_{o}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)-I_{n}\right) B_{\sigma}\left(d_{o}\right)\right]$.
We bound the contribution of the first term to $|L(u)|$; the second term can be treated similarly. Let $\tilde{L}(u)$ be the matrix when in (C.9) we replace $\nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right)$ by $\sum_{\sigma} c_{\sigma} \operatorname{tr}\left[B_{\sigma}\left(d_{o}\right)\right]$. Hence,

$$
\begin{equation*}
(\tilde{L}(u))_{l_{1}, l_{2}}=-\sum_{j=3}^{J} \frac{1}{(j-1)!} \sum_{\sigma \in \mathcal{S}_{j}} c_{\sigma} \sum_{l_{3}, \ldots, l_{j}=0}^{k} u_{l_{3}} \ldots u_{l_{j}} \frac{\operatorname{tr}\left[B_{\sigma}\left(d_{o}\right)\right]}{n}, \tag{C.26}
\end{equation*}
$$

$$
\begin{equation*}
\text { where } \frac{1}{n} \operatorname{tr}\left[B_{\sigma}\left(d_{o}\right)\right]=\frac{1}{2 \pi} \int_{-\pi}^{\pi} \cos \left(l_{1} x\right) \ldots \cos \left(l_{j} x\right) d x+E_{\sigma} \text {, } \tag{C.27}
\end{equation*}
$$

$E_{\sigma}$ being the approximation error. For each $\sigma$ and $j \geq 4$, the contribution of the integral in (C.27) to $(\tilde{L}(u))_{l_{1}, l_{2}}$ is $O\left(\int_{-\pi}^{\pi}\left|u^{t} \cos \right|{ }^{j-2}(x) d x\right)=O\left(\|u\|^{2}\right)$; hence its contribution to $|\tilde{L}(u)|$ is $k\|u\|^{2}=o\left(n^{-1 / 2-\delta} k\right)$. For $j=3$, we have

$$
\frac{1}{4 \pi} \sum_{l_{3}=1}^{k} u_{l_{3}} \int_{-\pi}^{\pi} \cos \left(l_{1} x\right) \cos \left(l_{2} x\right) \cos \left(l_{3} x\right) d x=\frac{1}{2}\left(u_{l_{1}+l_{2}} \mathbb{I}_{l_{3}=l_{1}+l_{2}}+u_{\mid l_{1}-l_{2}} \mathbb{I}_{l_{3}=\left|l_{1}-l_{2}\right|}\right),
$$

and the contribution of this term to $|\tilde{L}(u)|$ is of order $\sqrt{k}\|u\|=o\left(n^{-1 / 2+\epsilon} k\right)$.
Next we bound the contribution to $|\tilde{L}(u)|$ of the error term $E_{\sigma}$ in (C.27).
Note that we can write the last sum in (C.26) as
(C.28)

$$
\sum_{l_{3}, \ldots, l_{j}=0}^{k} u_{l_{3}} \ldots u_{l_{j}} \frac{\operatorname{tr}\left[B_{\sigma}\left(d_{o}\right)\right]}{n}=\frac{1}{n} \operatorname{tr}\left[\prod_{i=1}^{p} T_{n}\left(b_{i}(x) f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)\right]
$$

where

$$
\begin{equation*}
b_{i}(x)=\left(u^{t} \cos (x)\right)^{|\sigma(i)|-\delta_{1}(i)-\delta_{2}(i)} \cos \left(l_{1} \cdot\right)^{\delta_{1}(i)} \cos \left(l_{2} \cdot\right)^{\delta_{2}(i)} \tag{C.29}
\end{equation*}
$$

$\delta_{1}(i)=\mathbb{I}_{1 \in \sigma(i)}, \delta_{2}(i)=\mathbb{I}_{2 \in \sigma(i)}$ and $p=|\sigma|$. If $p \leq 3$, then Lemma 2.4 (supplement) implies that

$$
\begin{equation*}
E_{\sigma}=O\left((\sqrt{k}\|u\|)^{j-2} n^{\epsilon-1}\left[k^{2(3 / 2-\beta)_{+}}+k\right]\right)=o\left(n^{-1-\delta} k\right) \tag{C.30}
\end{equation*}
$$

If $p \geq 4$, then Lemma 2.6 (supplement) together with (C.28), with

$$
f=f_{d_{o}, k}, \quad f_{2 i}=b_{i} f_{d_{o}, k}, i \leq|\sigma|
$$

$L=k^{(3 / 2-\beta)_{+}}, M, m^{-1}=O(1), M^{(i)}=O\left((\sqrt{k}\|u\|)^{|\sigma(i)|-\delta_{1}(i)-\delta_{2}(i)}\right.$ and $L^{(i)}=O\left(k(\sqrt{k}\|u\|)^{|\sigma(i)|-\delta_{1}(i)-\delta_{2}(i)}\right)$, leads to the bound

$$
\begin{aligned}
& \sum_{l_{3}, \ldots, l_{j}=0}^{k} u_{l_{3}} \ldots u_{l_{j}} \frac{\operatorname{tr}\left[B_{\sigma}\left(d_{o}\right)\right]}{n}-\frac{1}{n} \operatorname{tr}\left[\prod_{i=1}^{|\sigma|} T_{n}\left(b_{i} f_{d_{o}, k}\right) T_{n}\left(\frac{1}{4 \pi^{2} f_{d_{o}, k}}\right)\right] \\
& =\mathbf{o}\left(k^{(3 / 4-\beta / 2)_{+}} n^{-1 / 2+\epsilon}\|u\|(\sqrt{k}\|u\|)^{j-3}\right)=\mathbf{o}\left(n^{-1 / 2-\delta}\right) .
\end{aligned}
$$

Using Lemma 2.1 (supplement) we finally obtain that

$$
\frac{1}{n} \operatorname{tr}\left[\prod_{i=1}^{|\sigma|} T_{n}\left(b_{i} f_{d_{o}, k}\right)\right]-\frac{1}{2 \pi} \int_{-\pi}^{\pi}\left(u^{t} \mathbf{c o s}\right)^{j-2}(x) \cos \left(l_{1} x\right) \cos \left(l_{2} x\right) d x=\mathbf{o}\left(n^{-1-\delta} k\right) .
$$

Therefore the contribution of the approximation error $E_{\sigma}$ in $|\tilde{L}(u)|$ is of order $\mathbf{o}\left(n^{-1 / 2-\delta} k\right)$. Using a similar argument we control the terms in the form $\operatorname{tr}\left[T_{n}\left(f_{o}\right)\left(T_{n}^{-1}\left(d_{o}, k\right)-I_{n}\right) B_{\sigma}\left(d_{o}\right)\right]$ and (C.21) is proved.

We now prove (C.22) and bound

$$
(G-\bar{G})_{l_{1}, l_{2}}=\sum_{j=2}^{J} \frac{1}{\bar{j}!} \sum_{\sigma \in \mathcal{S}_{j}}\left(c_{\sigma}-d_{\sigma}\right) \sum_{l_{3}, \ldots, l_{j}=0}^{k} r_{l_{1}, \ldots, l_{j}}^{(j)} u_{l_{3}} \ldots u_{l_{j}}
$$

with $r_{l_{1}, \ldots, l_{j}}^{(j)}$ as in (C.6). These are the approximation errors which occur when replacing $\frac{1}{n} \operatorname{tr}\left[T_{1, \sigma}\left(d_{o}, k\right)\right]$ and $\frac{1}{n} \operatorname{tr}\left[T_{2, \sigma}\left(d_{o}, k\right)\right]$ by their limiting integrals (see also (B.4)). Therefore, for each $\sigma \in \mathcal{S}_{j}, \sum_{l_{3}, \ldots, l_{j}=0}^{k} r_{l_{1}, \ldots, l_{j}}^{(j)} u_{l_{3}} \ldots u_{l_{j}}$ is a combination of terms of the form

$$
\frac{1}{n} \operatorname{tr}\left[\prod_{i=1}^{p} T_{n}\left(b_{i}(x) f_{d_{o}, k}\right) T_{n}^{-1}\left(f_{d_{o}, k}\right)\right]-\frac{1}{2 \pi} \int_{-\pi}^{\pi}\left(u^{t} \cos (x)\right)^{j-2} H_{k}(x) \cos \left(l_{1} x\right) \cos \left(l_{2} x\right) d x
$$

with $p \in\{|\sigma|,|\sigma|+1\}$ and the functions $b_{i}$ defined as in (C.29) apart from $b_{1}(x)=H_{k}(x)\left(\sum_{l=0}^{k} u_{l} \cos (l x)\right)^{|\sigma(1)|-\delta_{1}(1)-\delta_{2}(1)} \cos \left(l_{1} .\right)^{\delta_{1}(1)} \cos \left(l_{2} .\right)^{\delta_{2}(1)}$. Therefore, using the same construction as in (C.28)-(C.30), we obtain that

$$
\left|(G-\bar{G})_{l_{1}, l_{2}}\right|=O\left(n^{-1 / 2+\epsilon}\right), \quad|G-\bar{G}|=O\left(n^{-1 / 2+\epsilon} k\right)=o(1) .
$$

To prove (C.23), we use the just obtained bound on $|G-\bar{G}|$, and in addition establish a bound $|\bar{G}|$. We treat each term $G^{(j)}$ in $G(u)=\sum_{j=2}^{J} \frac{1}{j!} \sum_{\sigma \in \mathcal{S}_{j}}\left(c_{\sigma}-\right.$ $\left.d_{\sigma}\right) G^{(j)}(u)$ separately. First we show that $\left|\bar{G}^{(2)}\right|=O(1)$, which follows from definition (C.7), by which

$$
\operatorname{tr}\left[\left(\bar{G}^{(2)}\right)^{2}\right]=\sum_{l_{1}, l_{2}, l_{1}+l_{2} \geq k}^{k} \frac{1}{\left(l_{1}+l_{2}\right)^{2}} \leq 1
$$

Consequently, $\left|\bar{G}^{(2)}\right| \leq 1$. For $j \geq 3$, note that for all $0 \leq l_{1}, l_{2} \leq k$,

$$
\begin{aligned}
\left|\bar{G}_{l_{1}, l_{2}}^{(j)}(u)\right| & =\left|\sum_{l_{3}, \ldots, l_{j}=0}^{k} \gamma_{l_{1}, l_{2}, \ldots, l_{j}}^{(j)} u_{l_{3}} \ldots u_{l_{j}}\right| \\
& \leq \sum_{l_{4}, \ldots, l_{j}=0}^{k}\left|u_{l_{4}} \ldots u_{l_{j}}\right| \int_{-\pi}^{\pi}\left|H_{k}(x)\right|\left|\sum_{l_{3}=0}^{k} \cos \left(l_{3} x\right) u_{l_{3}}\right| d x \\
& \leq(\sqrt{k}\|u\|)^{j-3} \frac{\|u\|}{\sqrt{k}}=(\sqrt{k})^{j-4}(\|u\|)^{j-2} .
\end{aligned}
$$

Therefore, $\left|\bar{G}^{(j)}(u)\right| \leq k(\sqrt{k})^{j-4}(\|u\|)^{j-2}=(\sqrt{k}\|u\|)^{j-2}=o(1)$, for all $j \geq$ 3. Hence $|\bar{G}(u)|=\mathbf{O}(1)$, which combined with $\left\|\left(\tilde{I}_{k}^{+} L(u)\right)-1\right\|=\mathbf{O}_{P_{o}}(1)$ (see Lemma B.5) and (C.21)), imply that

$$
|\bar{D}(u)|=\left|\left(\tilde{I}_{k}+L(u)\right)^{-1} \bar{G}^{t}(u)\right|=O(1),
$$

uniformly over $\|u\| \leq 2 l_{0} \delta_{n}$. It follows that

$$
|D(u)| \leq\left\|\left(\tilde{I}_{k}+L(u)\right)^{-1}\right\|(|\bar{G}|+|G-\bar{G}|)=\mathbf{O}_{P_{o}}(1)
$$

This concludes the proof of (C.23); (C.24) directly follows from this result
since $\|u-\psi(u)\| \leq\left|d-d_{o}\right||D(u)|\|u\|$. Finally, we prove (C.25). We have

$$
\begin{aligned}
& \sum_{j=3}^{J} \frac{\left(v^{(j)}-u^{(j)}\right) \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}=-\left(d-d_{o}\right) \sum_{j=3}^{J} \sum_{l_{1}, \ldots, l_{j}=0}^{k}(D(u) u)_{l_{1}} u_{l_{2}} \ldots u_{l_{j}} \frac{\nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right)}{(j-1)!} \\
& \quad+\left(d-d_{o}\right)^{2} \sum_{j=3}^{J}\binom{j}{2} \sum_{l_{1}, \ldots, l_{j}=0}^{k}(D(u) u)_{l_{1}}(D(u) u)_{l_{2}} \ldots u_{l_{j}} \frac{\nabla_{l_{1}, \ldots, l_{j} l_{n}\left(l_{o}, k\right)}^{(j-1)!}}{} \\
& \quad+\ldots+(-1)^{J} \sum_{l_{1}, \ldots, l_{J}=0}^{k}(D(u) u)_{l_{1}}(D(u) u)_{l_{2}} \ldots(D(u) u)_{l_{J}} \frac{\nabla_{l_{1}, \ldots, l_{J}} l_{n}\left(d_{o}, k\right)}{(J-1)!}
\end{aligned}
$$

Using the same argument as in the proof of (C.21), we find that for all for all $j \geq 3$

$$
\begin{aligned}
& \sum_{l_{1}, \ldots, l_{j}=0}^{k}(D(u) u)_{l_{1}}(D(u) u)_{l_{2}} u_{l_{3}} \ldots u_{l_{j}} \frac{\nabla_{l_{1}, \ldots, l_{j}} l_{n}\left(d_{o}, k\right)}{j!} \\
& \left.=n \int_{-\pi}^{\pi}(D(u) u)^{t} \cos (x)\right)^{2}\left(u^{t} \cos (x)\right)^{j-2} d x \\
& \quad+(\sqrt{k}\|u\|)^{j-1} O\left(\sqrt{n} n^{\epsilon}(\sqrt{k}\|u\|)+k+\sqrt{n}\|u\| k^{1 / 2(3 / 2-\beta)_{+}}\right) \\
& =\mathbf{o}\left(n^{1-\delta} k^{-1}\right), \quad \text { for some } \delta>0 .
\end{aligned}
$$

Similarly, the higher-order terms in the above expression for $\sum_{j=3}^{J} \frac{\left(v^{(j)}-u^{(j)}\right) \nabla^{j} l_{n}\left(d_{o}, k\right)}{j!}$ can be shown to be $O\left(S_{n}(d)\right)$, which terminates the proof of Lemma C.1.

Biometris, Building 107
6708 PB Wageningen
The Netherlands
E-MAIL: willem.kruijer@wur.nl

3, Avenue Pierre Larousse
92245 Malakoff
France
E-mAIL: rousseau@ceremade.dauphine.fr


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