STATIONARY EXPONENTIAL FAMILIES

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A stationary exponential family is defined using transition densities which take the form of exponentiated symmetric k-linear forms on \mathbf{R}^d . Estimation is based on a mean value parametrization through a convex function on a finite-dimensional vector space. A consistency theorem and a central limit theorem are presented.

1. Introduction. We are concerned with defining and studying an exponential family for stationary sequences of random vectors X_1, X_2, \ldots in the finite-dimensional vector space $F = \mathbf{R}^d$. The main results are a consistency theorem and a central limit theorem (Theorems 2.1 and 2.2) for an estimator $\hat{\theta}$ of the parameter θ which indexes the process. We begin with some definitions and a preliminary discussion of the exponential family, and we present the main results in Section 2 and examples in Section 3.

Let E be the vector space of symmetric k-linear forms on F^k . The dimension of E is $\binom{k+d-1}{d-1}$, since a basis corresponds to a nonnegative integer solution of $x_1+\cdots+x_d=k$. Let μ be a reference probability measure on \mathbf{R}^d such that $Z_\theta=\int \exp(\theta(x_1,\ldots,x_k))\mu^k(dx_1,\ldots,dx_k)<\infty$ precisely on an open set $\Theta\subset E,\ \theta\in\Theta$. The Borel field on \mathbf{R}^d will be denoted B.

Construct a stationary process $\{X_i \in F, \ i \geq 1\}$ defined on the space F^{∞} as follows. Fix $\theta \in \Theta$ and define

$$\begin{split} Z(\,x_1,\ldots,x_k) &= \exp(\,\theta(\,x_1,\ldots,x_k)\,), \\ Z(\,x_1,\ldots,x_i) &= \int \! \exp(\,\theta(\,x_1,\ldots,x_k)\,) \,\mu^{k-i}(\,dx_k\,,dx_{k-1},\ldots,dx_{i+1})\,, \\ (1.1) &\qquad \qquad 1 \leq i \leq k\,, \\ Z &= \int \! \exp(\,\theta(\,x_1,\ldots,x_k)\,) \,\mu^k(\,dx_k\,,\ldots,dx_1)\,. \end{split}$$

The Z will appear with a subscript θ occasionally to avoid ambiguity. Let the transition density p (p_{θ} occasionally) be given by

$$p(x_1,...,x_k) = \frac{\exp(\theta(x_1,...,x_k))}{Z(x_1,...,x_{k-1})} = \frac{Z(x_1,...,x_k)}{Z(x_1,...,x_{k-1})}.$$

Let $\pi(x_1, \ldots, x_{k-1}, dx_k)$ denote the transition probability associated to the density p. The theorem of Ionescu Tulcea [see Ionescu Tulcea (1949) or

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Shiryayev (1984)] assures us a measure \mathbf{P}_{θ} on (F^{∞}, B^{∞}) which satisfies

$$\begin{split} \mathbf{P}_{\theta}(A_1 \times \cdots \times A_n \times F \times \cdots) &= \int_{A_1} \frac{Z(x_1)}{Z} \mu(dx_1) \cdots \\ &\times \int_{A_{k-1}} \frac{Z(x_1, \dots, x_{k-1})}{Z(x_1, \dots, x_{k-2})} \mu(dx_{k-1}) \cdots \\ &\times \int_{A_k} \pi(x_1, \dots, x_{k-1}, dx_k) \cdots \\ &\times \int_{A_n} \pi(x_{n-k+1}, \dots, x_{n-1}, dx_n) \end{split}$$

for every positive integer n and all sets $A_i \in B$, $1 \le i \le n$. The coordinate process X_1, X_2, \ldots on $(F^{\infty}, B^{\infty}, \mathbf{P}_{\theta})$ is a chain of finite order with transition probability π . One can verify that the process is strictly stationary and that the stationary measure on k coordinates for this process has density $\exp(\theta(x_1, \ldots, x_k))/Z$.

Define the sequence of random elements $\{Y_n\}$ in the dual space E' of the vector space E by

$$Y_n(\alpha) = \alpha(X_n, \ldots, X_{n+k-1})$$

for $\alpha \in E$. The vector Y_n is a k-fold tensor product defined via the symmetric k-linear forms.

Expectations and derivatives for the stationary exponential family of processes $\{\mathbf{P}_{\theta}\colon \theta \in \Theta\}$ are related as in the i.i.d. family, as we show in Proposition 1.1 below. For $\theta \in \Theta$, let \mathbf{E}_{θ} denote expectation with respect to the probability measure \mathbf{P}_{θ} on the measurable space F^{∞} . The symbol $\nabla \log Z$ denotes the derivative of the real-valued function Z, which is defined on E. Thus $\nabla \log Z \colon E \to E'$.

Proposition 1.1. $\mathbf{E}_{(\cdot)}(Y_n) = \nabla \log Z \colon \Theta \to E'$.

PROOF. For each $\alpha \in E$,

$$\begin{split} \mathbf{E}_{\theta}(Y_n)(\alpha) &= \int \alpha(x_1, \dots, x_k) \frac{Z_{\theta}(x_1, \dots, x_k)}{Z_{\theta}} \mu^k(dx_1, \dots, dx_k) \\ &= \nabla_{\!\theta} \log Z(\alpha), \end{split}$$

which proves that $\mathbf{E}_{\theta}(Y_n) = \nabla_{\!\!\theta} \log Z$. \square

In the following text we let $m_{\theta} = \mathbf{E}_{\theta}(Y_i) \in E'$. Consider now the problem of estimating the parameter θ . The ergodic theorem implies that the sequence $\overline{Y} = n^{-1} \sum^n Y_i \in E'$ has the property that $\overline{Y} \to \mathbf{E}_{\theta}(Y_i) = \nabla_{\theta} \log Z$. To simplify notation, let

$$f = \nabla \log Z \colon \Theta \to E'$$
.

Proposition 1.1 suggests that we define our estimator $\hat{\theta}$ for the parameter $\theta \in E$ by

$$\hat{\theta} = f^{-1}(\overline{Y})$$

since at least formally $\hat{\theta} = f^{-1}(\overline{Y}) \to f^{-1}(\mathbf{E}_{\theta}(Y_i)) = \theta$. If C is the closed convex hull of the support of the distribution of Y_i , then $\overline{Y} \in C$ and one needs to be sure that C is in the range of f and that f is bijective from E to C for (1.2) to make sense. These points will be addressed in the next section. We think of this estimation scheme very simply as choosing the parameter $\hat{\theta} \in E$ which makes the mean of Y_i under the law $\mathbf{P}_{\hat{\theta}}$ equal to the sample mean $\overline{Y} \in C$.

The results that follow in Section 2 justify the existence and continuity of f^{-1} , show asymptotic properties of \overline{Y} and translate these properties to $\hat{\theta}$ through f^{-1} .

- **2. Two limit theorems.** Our study of $\hat{\theta}$ below begins with regularity properties of f^{-1} to obtain the consistency result Theorem 2.1. Then we apply a central limit theorem of Rosenblatt (1971) to the multivariate Markov process (X_n,\ldots,X_{n+k-1}) in order to prove a central limit theorem for $\hat{\theta}$ at Theorem 2.2. Assume in what follows that:
- (2.1) (a) $\log Z_{\theta}$ is finite precisely for $\theta \in \Theta$, $\Theta \subset E$ open; (b) $Y_{i} \in \text{int } C$ a.s.

The process can be constructed with less than (2.1a), but we use all of (2.1) to prove results about $\hat{\theta}$.

LEMMA 2.1. Assume (2.1). Then the map $f = \nabla \log Z$: $\Theta \subset E \to \operatorname{int} C$ is bijective.

PROOF. The function $\log Z$ is strictly convex on Θ by Theorem 7.1 of Barndorff-Nielsen (1978) and (2.1b). Now we need only apply Theorems 5.33 and 9.2 of Barndorff-Nielsen (1978). \square

Lemma 2.1 makes sense of the estimation procedure (1.2), since now we have with assumption (2.1b) that $\hat{\theta}$ exists and is well defined with probability 1. Next we look for further regularity properties of f.

Lemma 2.2. Assume (2.1). Then $f = \nabla \log Z$ has differentiable inverse f^{-1} : int $C \to \Theta \subset E$.

PROOF. By Lemma 2.1, f^{-1} exists. To show it is differentiable, it is enough to show that f is differentiable without critical points. However, $D_{\theta}f$ is a linear map from E to E' corresponding in the standard way to a bilinear form on $E \times E$ which satisfies

$$D_{\theta}f(\alpha,\alpha) = \mathbf{E}_{\theta}(\langle \alpha, Y_1 - m_{\theta} \rangle^2)$$

for any $\alpha \in E$. $D_{\theta}f$ is, in fact, positive definite, since (2.1b) implies that the law of Y_i is not concentrated on an affine subspace of E'. Therefore, the matrix for $D_{\theta}f$ is invertible, and so f has no critical points. \square

Theorem 2.1. Assume (2.1). Then the estimator $\hat{\theta} \in E$ converges a.s. to θ .

PROOF. The map f^{-1} is continuous, so the formal argument following (1.2) is valid. \Box

We will prove below a central limit theorem for $\hat{\theta}$ from results of Rosenblatt (1971) for Markov chains. We start with the following observation.

Lemma 2.3. Assume (2.1). The process $\{(X_n,\ldots,X_{n+k-1}):\ n\geq 1\}$ is a Markov chain on F^k with transition probabilities π_k given by

$$\begin{split} \pi_k \big((x_0, \dots, x_{k-1}), (A_1, \dots, A_k) \big) \\ &= \delta_{(x_1, \dots, x_{k-1})} (A_1 \times \dots \times A_{k-1}) \pi(x_1, \dots, x_{k-1}, A_k). \end{split}$$

Let $L^{2,\,k}$ denote the set of real-valued measurable functions on the product space F^k such that

$$||g||_2^2 = \mathbf{E}_{\theta}(g^2(X_1,\ldots,X_k)) < \infty.$$

This is the L^2 norm for the stationary probability measure on k-coordinates having density $Z(x_1, \ldots, x_k)/Z$. Define the operator T on $L^{2, k}$ by

$$(Tg)(x_1,...,x_k) = \int \pi_{\theta}(x_2,...,x_k,dx_{k+1})g(x_2,...,x_{k+1}).$$

We will say that $g \perp 1$ if $\mathbf{E}_{\theta}(g(X_1, \ldots, X_k)) = 0$. Consider the following L^2 norm condition:

(2.2)
$$\sup_{g \perp 1} \frac{\|T^n g\|_2}{\|g\|_2} \to 0.$$

Lemma 2.4. Assume (2.1) and (2.2). Then the sequence $\sum_{1}^{n}(Y_{i}-m_{\theta})/\sqrt{n}$ converges in distribution to the multivariate normal law $\mathbf{N}(0,A)$, where A is the symmetric bilinear form on $E\times E$ defined by

$$A(\alpha, \alpha) = \lim \frac{1}{n} \mathbf{E}_{\theta} \left(\left(\sum_{i=1}^{n} \langle \alpha, Y_{i} - m_{\theta} \rangle \right)^{2} \right).$$

PROOF. Consider the usual operator for the multivariate Markov chain with the kernel π_k . Then (2.2) is the L^2 norm condition of Rosenblatt [(1971), page 206]. Apply Theorem 2 of Rosenblatt [(1971), page 217], to the sequence $S_n = \langle \alpha, \Sigma_1^n(Y_i - m_\theta) \rangle$ [where $f_{i,n} = (1/\sqrt{n})\alpha(X_i, \dots, X_{i+k-1})$ and $k_n = n$ in

his notation] to see that $S_n/\sqrt{n}\to N(0,A(\alpha,\alpha))$. Then it is also true that for each $\alpha,\beta\in E$ the quantity

$$rac{1}{n}\mathbf{E}_{ heta}igg(\sum_{1}^{n}ra{lpha},Y_{i}-m_{ heta}ig)\sum_{1}^{n}raket{eta,Y_{i}-m_{ heta}}$$

converges to a number, say $A(\alpha, \beta) < \infty$, and A is necessarily symmetric and bilinear. Thus the covariance matrix A is of the stated form.

However, this implies that the sequence of vectors $\sum_{i=1}^{n} (Y_i - m_{\theta}) / \sqrt{n} \in E'$ converges in distribution to $\mathbf{N}(0, A)$. \square

Consider the condition (2.2). If there exists $g \in L^1(\mu)$ such that

$$(2.3) p(x_1, ..., x_{k-1}, \cdot) \leq g(\cdot),$$

then (2.2) and the above central limit theorem hold. The condition (2.3) gives uniform integrability of the kernels π_k , which is stronger than the Doeblin condition and has been used extensively for central limit theorems [Doob (1953)]. Rosenblatt [(1971), Theorem 1, page 211] shows that the Doeblin condition implies the L^2 norm condition, and generally (2.3) is easier to check. Condition (2.2) means in a precise way that the multivariate chain (X_n, \ldots, X_{n+k-1}) is asymptotically uncorrelated [Rosenblatt (1971), page 207].

For k=2, the process (X_n) is Markov with ordinary transition operator π . Now π is self-adjoint on the space L^2 of functions on F with the stationary distribution, since

$$\langle g, \pi h \rangle = \int \mu(dx_1) Z(x_1) / Z g(x_1) \int \mu(dx_2) Z(x_1, x_2) / Z(x_1) h(x_2)$$

$$= \int \mu(dx_2) Z(x_2) / Z h(x_2) \int \mu(dx_1) Z(x_1, x_2) / Z(x_2) g(x_1)$$

$$= \langle \pi g, h \rangle$$

using the symmetry of $Z(\cdot,\cdot)$. Note that $T^{n+1}g(x_1,x_2)=\pi^nTg(x_2)$, and if $g\perp 1$, then h=Tg is orthogonal to 1 in L^2 and thus (2.2) is satisfied provided

$$\sup_{h \perp 1} \frac{\|\pi^n h\|_2}{\|h\|_2} \to 0.$$

In particular, if π has a complete set of eigenfunctions (necessarily orthogonal), then for $h \perp 1$,

$$\|\pi^n h\|_2 \leq |\lambda|^n \|h\|_2,$$

where λ is the second largest eigenvalue of π in absolute value, or the largest eigenvalue of π when restricted to the orthogonal complement of the function 1. The condition (2.2) is then satisfied if $|\lambda| < 1$, which will be the case in Examples 3.3 and 3.5.

Theorem 2.2 below is the central limit theorem for $\hat{\theta}$. Recall that as a bilinear map on $E \times E$, $D_{\theta} f$ satisfies $D_{\theta} f(\alpha, \alpha) = \mathbf{E}_{\theta} (\langle \alpha, Y_1 - m_{\theta} \rangle^2)$ and can thus be interpreted as the covariance matrix for the vector Y_1 .

THEOREM 2.2. Assume (2.1) and (2.2). Then the vector $\sqrt{n}(\hat{\theta} - \theta)$ converges in distribution to $\mathbf{N}(0, B)$, where B is the symmetric bilinear form on $E' \times E'$ defined by

$$B(v,v) = A([D_{\theta}f]^{-1}(v), [D_{\theta}f]^{-1}(v)).$$

PROOF. Since $\hat{\theta} = f^{-1}(\overline{Y})$ and f^{-1} is differentiable at $f(\theta)$, it follows that $\hat{\theta} - \theta = D_{f(\theta)} f^{-1}(\overline{Y} - f(\theta)) + o(\overline{Y} - f(\theta)),$

where $D_{f(\theta)}f^{-1}$ is a linear map from E' to E and $o(\cdot)$ is a function from E' to E such that $o(x)/\|x\| \to 0$ as $\|x\| \to 0$. We can use the standard argument for transferring central limit theorems through differentiable maps to see that for each $v \in E'$,

$$\langle \sqrt{n} (\hat{\theta} - \theta), v \rangle \rightarrow N(0, A([D_{\theta}f]^{-1}(v), [D_{\theta}f]^{-1}(v))).$$

However, this proves the multivariate assertion that $\sqrt{n}(\hat{\theta} - \theta)$ converges in distribution to N(0, B). \square

3. Examples.

EXAMPLE 3.1. Let d=1. Then E and E' have dimension $\binom{k+d-1}{d-1}=1$, regardless of k. The parameter θ is simply an element of \mathbf{R}^1 , and the process Y_n takes values in \mathbf{R}^1 as well as is given by $Y_n=\prod_{0}^{k-1}X_{n+i}$. Thus

$$m_{ heta} = \int_{\mathbf{R}^k} x_1 imes \cdots imes x_k rac{\exp(\, heta x_1 imes \cdots imes x_k)}{Z_o} \, d\mu^k(\,x_1, \ldots, \,x_k).$$

If condition (2.2) holds, let $\sigma^2 = \lim(1/n)\mathbf{E}_{\theta}(\Sigma^n(Y_i - m_{\theta}))^2$. Then Σ^nY_i/\sqrt{n} converges in distribution to $N(m_{\theta}, \sigma^2)$. Since $[D_{\theta}f]^{-1} = 1/\mathbf{E}_{\theta}(Y_1 - m_{\theta})^2$, it follows that $\sqrt{n}(\hat{\theta} - \theta)$ converges in distribution to $N(0, \sigma^2/[\mathbf{E}_{\theta}(Y_1 - m_{\theta})^2]^2)$.

EXAMPLE 3.2. The special case where d=1 and k=2 corresponds to a Markov chain on the real line. The stationary density is given by Z(x)/Z. The chain is reversible since Z(x)p(x,y)/Z=Z(y)p(y,x)/Z, which leads to a self-adjoint transition operator π , but not symmetric.

Let the reference measure μ be Lebesgue measure on (0, 1). Then we can take $\Theta = \mathbf{R} = E$ and the transition density p(x, y) is given by

$$p(x,y) = \frac{\theta x e^{\theta x y}}{e^{\theta x} - 1},$$

which is bounded uniformly in x on (0,1) and satisfies conditions (2.1) and (2.2) (in fact, the chain satisfies the Doeblin condition with uniformly μ -integrable transition functions). The map f is given by

$$f(\theta) = \int_0^1 \int_0^1 xy e^{\theta xy} dx dy / \int_0^1 \int_0^1 e^{\theta xy} dx dy$$
$$= \frac{1}{\theta} \left(\frac{e^{\theta} - 1}{c_{\theta}} - 1 \right),$$

where $c_{\theta} = \int_0^{\theta} (e^x - 1)/x \, dx$. The function f is monotone increasing from the real line $\Theta = (-\infty, \infty)$ onto the interval $(0,1) = \operatorname{int}(C)$ (cf. Lemma 2.1). Now $\overline{Y} = \sum_{i=1}^{n} X_i X_{i+1}/n$ and the estimator $\hat{\theta}$ is the unique solution to

$$f(\hat{\theta}) = \overline{Y} \in (0,1).$$

EXAMPLE 3.3. Suppose we are interested in a standard Gaussian reference measure. Then the transition density p(x, y) with respect to Lebesgue measure on **R** becomes

$$p(x,y) = \frac{1}{\sqrt{2\pi}} \exp\left(-\frac{(y-\theta x)^2}{2}\right),$$

and when $\theta \in (-1,1)$ this represents the transition probability for a process in discrete time analogous to the Ornstein–Uhlenbeck process. The sequence X_0, X_1, \ldots can also be represented as an autoregressive series with $X_n = \xi_n + \theta X_{n-1} = \sum_0^\infty \theta^i \xi_{n-i}, \ \xi_i \sim \text{i.i.d.} \ N(0,1).$ Set $\Theta = (-1,1)$. Conditions (2.1) are satisfied and the invariant marginal

Set $\Theta=(-1,1)$. Conditions (2.1) are satisfied and the invariant marginal distribution on X_n is $N(0,1/(1-\theta^2))$. Condition (2.2) is also satisfied, since the self-adjoint operator π on L^2 takes the form

$$\pi h(x) = \int h(y) \frac{\exp(-(y-\theta x)^2/2)}{\sqrt{2\pi}} dy.$$

One can check that the spectrum of π is $1, \theta, \theta^2, \ldots$ and that these numbers correspond to Hermite polynomials for the stationary distribution $N(0, 1/(1-\theta^2))$. By the remarks following Lemma 2.4, the L^2 norm condition is satisfied.

The map $f: \Theta \to \mathbf{R}$ is given by

$$f(\theta) = \frac{\theta}{1 - \theta^2},$$

which is an increasing map from $\Theta = (-1,1)$ onto $E = (-\infty,\infty)$, with asymptotes at ± 1 . This is of course the covariance of X_n and X_{n+1} with joint bivariate normal distribution having mean 0 and covariance matrix R given by

$$R^{-1} = \begin{pmatrix} 1 & -\theta \\ -\theta & 1 \end{pmatrix}.$$

One can solve explicitly $f(\hat{\theta}) = \overline{Y}$ to find the estimator for θ . With some straightforward calculation it is seen that $Cov(Y_i, Y_j) = \theta^2 Cov(Y_i, Y_{i-1})$ for

 $j \ge i + 2$, and

$$\sigma^2 = \lim rac{1}{n} \mathrm{Var}_{ heta} \Big(\sum^n Y_i \Big) = rac{1 - heta^4 + 4 heta^2}{\left(1 - heta^2
ight)^3},$$
 $\mathrm{Var}(Y_i) = f'(heta) = rac{1 + heta^2}{\left(1 - heta^2
ight)^2},$
 $\sqrt{n} \left(\hat{ heta} - heta
ight)
ightarrow N igg(0, 1 + heta^2 - rac{8 heta^4}{\left(1 + heta^2
ight)^2} igg).$

The information $I_n(\theta)$ in the sample $\{X_0, X_1, \dots, X_n\}$, given by

$$I_n(\, heta\,) = \mathbf{E}_{ heta} \left[\left(\partial_{ heta} \log(1-\, heta^{\,2}) \, - X_0^2(1-\, heta^{\,2}) \, - \, \sum_1^n \left(\,X_i \, - \, heta X_{i-1}
ight)^2
ight]^2
ight]$$

satisfies $I_n(\theta)/n \to 1/(1-\theta^2)$. The maximum likelihood estimator $\hat{\theta}_{\rm ML}$ is the value of θ which maximizes the quantity $\log(1-\theta^2)-X_0^2(1-\theta^2)-\sum_1^n(X_i-\theta X_{i-1})^2$. It is known that $\hat{\theta}_{\rm ML}$ satisfies $\sqrt{n}\,(\hat{\theta}_{\rm ML}-\theta)\to N(0,1-\theta^2)$ [Box and Jenkins (1970), pages 280–281] and that this asymptotic variance is also attained by the sample correlation r_1 . Thus the asymptotic variance for $\hat{\theta}$ in this example is slightly greater than the variance of these estimators.

EXAMPLE 3.4. Let d=2 and let k=2. Then E and E' have dimension 3, and E consists of symmetric bilinear forms on $\mathbf{R}^2 \times \mathbf{R}^2$. Elements of E can be identified with symmetric 2×2 matrices and a basis consists of the three vectors

$$(e_1,e_2,e_3)=\left(\begin{pmatrix}1&0\\0&0\end{pmatrix},\begin{pmatrix}0&1/\sqrt{2}\\1/\sqrt{2}&0\end{pmatrix},\begin{pmatrix}0&0\\0&1\end{pmatrix}\right),$$

which are orthogonal for the entrywise inner product. E is then a three-dimensional subspace of the vector space of all 2×2 matrices. A corresponding dual basis (e'_1, e'_2, e'_3) for E' has the same matrix representation and $\langle e_i, e'_j \rangle = \delta_{ij}$ for the entrywise inner product. For example, if we think of \mathbf{x}, \mathbf{y} in the state space $F = \mathbf{R}^2$ as column vectors, then

$$e_1(\mathbf{x},\mathbf{y}) = \mathbf{x}^T \begin{pmatrix} 1 & 0 \\ 0 & 0 \end{pmatrix} \mathbf{y}.$$

With this identification, $f(\theta)$ is an element of E' determined by coordinates

Also $D_{\theta}f$ is a linear map from E to E' given by a 3×3 matrix (d_{ij}) with columns $D_{\theta}f(e_i)$, the entries of which are determined by the effect of $D_{\theta}f(e_i)$ on each basis element:

$$\begin{aligned} d_{ij} &= \big[D_{\theta}f(e_i)\big](e_j) = \mathbf{E}_{\theta}\big[\big(e_i(X_1, X_2) - \mathbf{E}_{\theta}\big[e_i(X_1, X_2)\big]\big) \\ &\times \big(e_j(X_1X_2) - \mathbf{E}_{\theta}\big[e_j(X_1, X_2)\big]\big)\big]. \end{aligned}$$

The bilinear form A on $E \times E$ can also be represented as a 3×3 symmetric matrix (a_{ij}) , where

$$a_{ij} = \lim \frac{1}{n} \mathbf{E}_{\theta} \left(\sum_{m=1}^{n} \langle e_i, Y_m - m_{\theta} \rangle \sum_{m=1}^{n} \langle e_j, Y_m - m_{\theta} \rangle \right).$$

EXAMPLE 3.5. Take d=2 and k=2 as above and let the reference measure μ on \mathbf{R}^2 be the product of two independent N(0,1) distributions. Let $\Theta \subset E$ consist of those symmetric matrices

$$heta = \left(egin{matrix} a & b \\ b & c \end{array}
ight)$$

such that the matrix

$$M_{ heta} = egin{pmatrix} 1 - a^2 - b^2 & -b(a+c) \ -b(a+c) & 1 - b^2 - c^2 \end{pmatrix} = id - heta^2$$

is positive definite. This means that the eigenvalues of θ are in (-1,1). One can easily show that conditions (2.1) are satisfied. One can show that condition (2.2) is satisfied, as we did in Example 3.3, by showing that the self-adjoint operator π on L^2 has largest eigenvalue equal to the largest eigenvalue θ_1 of the matrix θ , when π is restricted to the subspace of L^2 orthogonal to the function 1.

One finds that $Z_{\theta}^2 = 1/\det M_{\theta}$, and the transition density $p(\mathbf{x}, \mathbf{y})$ with respect to Lebesgue measure takes a form analogous to the one-dimensional situation,

$$p(\mathbf{x}, \mathbf{y}) = \frac{1}{2\pi} \exp\left(\frac{-\|\mathbf{y} - \theta \mathbf{x}\|^2}{2}\right),$$

so the step has the Gaussian distribution with mean $\theta \mathbf{x}$. The stationary marginal distribution on $F = \mathbf{R}^2$ is bivariate normal $N(0, M_{\theta}^{-1})$ and the stationary distribution on $F \times F = \{(x_1, x_2, y_1, y_2)\}$ is Gaussian with mean 0 and covariance matrix R such that

$$R^{-1} = \begin{pmatrix} 1 & 0 & -a & -b \\ 0 & 1 & -b & -c \\ -a & -b & 1 & 0 \\ -b & -c & 0 & 1 \end{pmatrix}.$$

In particular,

$$egin{aligned} f(heta)(e_1) &= rac{1}{Z_ heta} \int x_1 y_1 \exp(\mathbf{x}^T\! heta \mathbf{y}) \mu imes \mu(d\mathbf{x}, d\mathbf{y}) = rac{a + cb^2 - ac^2}{\det M_ heta}, \ f(heta)(e_2) &= rac{1}{Z_ heta} \int rac{y_1 x_2 + x_1 y_2}{\sqrt{2}} \exp(\mathbf{x}^T\! heta \mathbf{y}) \mu imes \mu(d\mathbf{x}, d\mathbf{y}) = rac{\sqrt{2} \, b(1 + ac - b^2)}{\det M_ heta}, \ f(heta)(e_3) &= rac{1}{Z_ heta} \int x_2 y_2 \exp(\mathbf{x}^T\! heta \mathbf{y}) \mu imes \mu(d\mathbf{x}, d\mathbf{y}) = rac{1 - a^2 - b^2}{\det M_ heta}. \end{aligned}$$

Thus we can represent $f(\theta) \in E'$ as the symmetric 2×2 matrix

$$f(\, heta) = rac{1}{\det M_{ heta}} igg(egin{array}{ccc} a + cb^2 - ac^2 & b(1 + ac - b^2) \ b(1 + ac - b^2) & 1 - a^2 - b^2 \end{array} igg) \ = heta(id - heta^2)^{-1}.$$

Finally, the statistic $\overline{Y} \in E'$ takes the following form. Let $X_i = (X_{i,1}, X_{i,2}) \in F = \mathbf{R}^2$. Then

$$\overline{Y}(e_i) = e_i \frac{1}{n} \sum_{1}^{n} (X_i, X_{i+1}),$$

and this gives the formulas

$$egin{aligned} & \overline{Y}(e_1) = rac{1}{n^2} \Sigma X_{i,1} \Sigma X_{i+1,1}, \ & \overline{Y}(e_2) = rac{1}{\sqrt{2} \, n^2} ig[\, \Sigma X_{i,2} \Sigma X_{i+1,1} + \Sigma X_{i,1} \Sigma X_{i+1,2} ig], \ & \overline{Y}(e_3) = rac{1}{n^2} \Sigma X_{i,2} \Sigma X_{i+1,2}, \end{aligned}$$

which in matrix form becomes

$$\overline{Y} = egin{pmatrix} \overline{Y}(e_1) & \overline{Y}(e_2)/\sqrt{2} \\ \overline{Y}(e_2)/\sqrt{2} & \overline{Y}(e_3) \end{pmatrix}.$$

We remark that the exponential family presented here is a restrictive parametric model for a stationary process that is different from typical time series models. The advantage of this is a very simple estimation scheme which is based on studying a sample mean in a finite-dimensional vector space. The sample mean itself is an easy candidate for an ergodic theorem and a central limit theorem. Then a smooth map transfers its statistical properties to the estimator $\hat{\theta}$.

One can ask about the efficiency of the estimator $\hat{\theta}$ which we have introduced based on natural convexity considerations. Example 3.3 indicates that $\hat{\theta}$ does not always have minimum asymptotic variance, although it is close. In the i.i.d. model (k=1), $\hat{\theta}$ is of course the MLE and an optimal large deviation property for the estimator was recently proved by Kester and Kallenberg (1986). A similar result for $k \geq 2$ will require first a careful look at the large deviation properties of \overline{Y} .

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