# ON THE CRAMER-RAO INEQUALITY IN AN INFINITE DIMENSIONAL SPACE

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ABSTRACT. An infinite-dimensional analogue of the Cramer-Rao inequality is given. The technique of smooth measures are used. Conditions of regularity are given and under these conditions a variant of maximal likelihood principle for the infinite-dimensional case is proposed. The consistency property of the maximum likelihood estimate is given.

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The Cramer-Rao (C-R) inequality and ensuing from it consequences play fundamental role in statistical analysis. Many important problems are being solved on the grounds of that analysis. But a range of such problems do not involve situations which are connected with random processes. Therefore the key point is to extend the methods of the C-R inequality to an infinite dimensional case. Contemporary state of the infinite dimensional analysis allows one to consider many basic problems of statistics from a more general point of view.

The approach based on the analysis of sensibility of a family of probabilistic measures is well-known (see, for e.g., E. Pitmen). The theory of smooth measures ([2]–[3]) gives us a good chance for generalization in this direction. The present paper realizes this chance by an example in which the C-R inequality is generalized to an infinite dimensional case. The outlines of such ideas have actually been given by E. Gobet in [4]. The idea consists in application of the theory of P. Malliavin calculus ([5]–[8]). In their work [9], J.M. Corcuera and A. Kohatsu-Higa have used the technique of stochastic calculus of variations (Malliavin calculus) and obtained the

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results for a finite dimensional (more precisely, for a one-dimensional) case. In the present work the use will be made of the methods of smooth measures allowing us to formulate a general look at the questions dealt with the C-R inequality and related with them problems. In [10], the theory of smooth measures has been used to estimate a logarithmic measure derivative.

## 1. The Logarithmic Measure Derivative

Let  $\{\Omega, \Im, P\}$  be a complete probabilite space.  $X = X(\omega; \theta)$  is a randon element with values in a linear space E and parameter  $\theta \in \Theta$ , where  $\Theta$  a subset of a separable real Banach space  $\Xi$  has the norm  $\|\cdot\|_{\Xi}$ .

The basic problem of statistics is to estimate an unknown parameter  $\theta$ . This estimation should be based on observations for realization of the given randon value  $X_1, X_2, \ldots, X_n, \ldots$ , a sequence (a sampling) of independent and identical to X distributed random values. We are required to construct such a statistics  $T = T(X_1, X_2, \ldots, X_n)$  which will be optimal (in a sense motivated in advance) to estimate  $\theta$ .

Usually, this situation generates a sequence of statistical structures  $\{\aleph, \Re, (P(\theta; \cdot), \theta \in \Theta)\}$ , where  $\aleph = E^n (n = 1, 2, ..., \infty)$  is a linear space generated by a sequence of random values  $X_1, X_2, ..., X_n$ ,  $\Re$   $\sigma$  is the algebra generated by observable sets, and  $\{P(\theta; \cdot), \theta \in \Theta\}$  is a system of probabilite measures (distributions) generated by the vector  $Y = (X_1, X_2, ..., X_n)$  by virtue of the relation  $P(\theta; A) = P(Y^{-1}(A))$ ,  $A \in \Re$ . In the classical statistics, the main object of investigation is, namely, that statistical structure  $\{\aleph, \Re, (P(\theta; \cdot), \theta \in \Theta)\}$ .

On the other hand, there exist a vast variety of problems in which  $X = X(\omega; \theta)$  is just the function, convenient to operate with by imposing certain analytic requirements of smoothness (the conditions of regularity) with respect to the parameter  $\theta$ . In addition, there arises a good possibility to apply an apparatus of stochastic calculus of variations.

Thus we obtain double calculus: the first one relies on the study of properties of the statistical structure  $\{\aleph, \Re, (P(\theta; \cdot), \theta \in \Theta)\}$ , the smoothness of the measure family  $P(\theta, \cdot)$ , and the second, the direct stochastic methods whose object of investigation is  $X(\omega, \theta)$ .

In this connection, we will be interested in the family of distributions  $\{P(\theta; A), \theta \in \Theta, A \in \Re\}$  from the point of view of their smoothness with respect to the both parameters  $\theta$  and A. Here we cite some definitions, therminological agreements and properties.

Throughout the paper, it will be assumed that  $\aleph$  is a separable, reflexive Banach space. For every fixed  $\theta \in \Theta$ ,  $P(\theta, \cdot)$  is of positive measure. If  $h \in \aleph$  is some vector, then by  $P_h(\theta; A)$  we denote a measure obtained by the shift  $P_h(\theta; A) = P(\theta; A + h)$ .

We say that the measure  $P(\theta;\cdot)$  is differentiable along the vector h, if there exists a bounded linear functional on  $\aleph$ , denoted by  $d_h P(\theta;\cdot)$ , such that for every  $A \in \Re$ , the equality

$$P_h(\theta; A) - P(\theta; A) = d_h(\theta; A)h + \alpha(\theta, A; h),$$

holds, where  $\alpha(\theta, A; h)$  is the function, such that  $\alpha(\theta, A; th) = o(t)$ ,  $t \in R$ . This is the so-called measure differentiability due to Fomin (the detailed account of another, somewhat different from Fomin's differentiability, definitions, properties and interconnections between them can be found in monograph [2]).

In the case, where  $\aleph$  is a separable, real Hilbert space with scalar product  $(\cdot,\cdot)_{\aleph}$  and norm  $||h||_{\aleph}$ ,  $h \in \aleph$ , we write  $P_h(\theta;A) - P(\theta;A) = (d_hP(\theta;A),h)_{\aleph} + \alpha(\theta,A;h)$ , and sometimes (when this does not put us to confusion) under the derivative  $d_hP(\theta;\cdot)$  will be understood an element of Hilbert space. Surely, the function  $d_hP(\theta;\cdot)$  is the  $\sigma$ -additive (of alternating signs) measure on  $\Re$ .

The measure derivative of higher order is defined by iteration in the course of determination of the derivative. Thus, for example,  $d_k d_h P(\theta; \cdot) = d_k (d_h P(\theta; \cdot)h)k$ ,  $k, h \in \aleph$ . In particular, in the case of Hilbert space  $\aleph$ , we have  $(d_{h,h}^{(2)} P(\theta; \cdot)h, h)_{\aleph} = (d_h (d_h P(\theta; \cdot), h)_{\aleph}, h)_{\aleph}$ .

For the n times differentiable measure, the expansion

$$P(\theta; A + h) - P(\theta; A) = \sum_{k=1}^{n} \frac{1}{k!} d_{h,\dots,h}^{(k)} P(\theta; A) + \alpha(\theta, A; h),$$

is valid, where  $\varphi(t) = \alpha(\theta, A; th)$  is the (n-1)-multiply differentiable function, vanishing together with its derivatives in zero faster than the corresponding powers t.

The function  $\psi_{\theta}(t) = P(\theta, A + th)$  is nonnegative and everywhere differentiable. If  $P(\theta; A) = 0$ ,  $A \in \Re$ , then the point t = 0 will turn out to be the point of minimum for the function  $\psi_{\theta}(t)$ . Therefore  $d_h(\theta; A) = 0$ . Hence, by the Radon-Nikodym's theorem, there exists a measurable function  $\beta_{\theta}(x; h)$  such that  $\frac{d_h P(\theta; dx)}{P(\theta; dx)} = \beta_{\theta}(x, h)$ . This function is called logarithmic measure derivative along the vector  $h \in \Re$ . The logarithmic derivative  $\beta_{\theta}(x, h)$  is linear in the second argument. The vector h is called an admissible direction for measure  $P(\theta; \cdot)$ . A set of all admissible directions is called an admissible subspace.

**Example 1.** Let  $H_+ \subset H \subset H_-$  be a triple of Hilbert spaces whose embedding operator  $i: H_+ \to H$  is the Hilbert-Schmidt's operator. Such a triple is called a Hilbert-Schmidt structure, or an equipped Hilbert space with quasi-kernel embeddings. Let  $\gamma_{\theta}$  be Gaussian measure in  $H_-$  with a unit correlation operator in H with  $\theta$  mean,  $\theta \in H_-$ . An admissible space

for  $\gamma_{\theta}$  is  $H_{+}$ . In addition, if  $h \in H_{+}$ , then the logarithmic derivative of measure  $\gamma_{\theta}$  along h is  $(\theta - x, h)_{H}$ .

In the theory of differentiable measures, of great importance is the fact that the formula of integration by parts is valid. Let  $\aleph$  be the separable, real Hilbert space and f(x) be a functional in that space. Assume that there exists its derivative directed with respect to the vector  $h \in \aleph$ , and  $d_h f(x) = \lim_{t\to 0} t^{-1}[f(x+th) - f(x)]$ , and  $d_h f(\cdot) \in L_1(P(\theta;\cdot))$  for the fixed  $\theta \in \Theta$ . Then, if measure  $P(\theta;\cdot)$  is differentiable with respect to the direction h, then (see [2])

$$\int_{\aleph} (d_h f(x), h)_{\aleph} P(\theta; dx) = -\int_{\aleph} f(x) d_h P(\theta; dx) =$$

$$= -\int_{\aleph} f(x) \beta_{\theta}(x; h) P(\theta; dx). \tag{1}$$

We can define logarithmic derivative along nonconstant directions (the so-called logarithmic gradient). Equality (1) may serve as a basis for such a definition, or we can act analogously to what we have done in determining the measure derivative along constant directions.

Let  $z(x): \aleph \to \aleph$  be the differentiable vector field possessing bounded derivative  $\sup_{x \in \aleph} ||z'(x)|| < \infty$ . An integral flow, corresponding to z(x), we denote by  $S_t$ ,  $t \in R$ . This implies that

$$\frac{dS_t}{dt} = z(S_t), \quad S_0 = I.$$

The family of measures  $(P(\theta;\cdot)\theta \in \Theta)$  is associated with a class of measures  $(P_t(\theta;\cdot)\theta \in \Theta, t \in R)$  due to the transformation  $P_t(\theta;A) = P(\theta;S_t^{-1}(A)), A \in \Re$ .

We say that the measure  $P(\theta;\cdot)$  is differentiable along the vector field z(x), if there exists the measure (necessarily of alternating signs)  $D_z P(\theta;A)$  such that for any bounded and differentiable function  $\varphi: \aleph \to R, \ \varphi \in C^1(\aleph;R)$  we have

$$\int_{\mathbb{N}} \varphi(x) D_z P(\theta; dx) = -\lim_{t \to 0} \int_{\mathbb{N}} \varphi(x) \frac{P_t(\theta) - P(\theta)}{t} (dx).$$

Hence after the transformation, we obtain

$$\int_{\aleph} \varphi(x) D_z P(\theta; dx) = -\int_{\aleph} \varphi'(x) z(x) P(\theta; dx).$$

If, in addition,  $D_z P(\theta; \cdot) \ll P(\theta; \cdot)$ , then the Radon-Nikodym's density is called the logarithmic derivative  $P(\theta; \cdot)$  along the vector field z(x):

$$\beta_{\theta}(x;z) = \frac{D_z P(\theta; dx)}{P(\theta; dx)}.$$

Let H be the embedding in the  $\aleph$  Hilbert space whose embedding operator is the Hilbert-Schmitd's operator. Then we can consider the Hilbert-Schmidt's structure  $\aleph^* \subset H \subset \aleph$ . We distinguish an important class of measures  $\mathcal{L}$  for which there exists a measurable, locally bounded function  $\lambda : \aleph \to \aleph$  such that for every constant direction  $h \in \aleph^*$  there exists the logarithmic derivative along h of the form  $\beta_{\theta}(x;h) = \lambda(\theta;x)h = (\lambda(\theta,x),h)_H$ . In this case we say that the measure possesses the logarithmic gradient  $\lambda(\theta;x)$ .

If  $P(\theta) \in \mathcal{L}$ , and the vector field  $z : \aleph \to \aleph^*$  is bounded together with its derivative, then for the measure  $P(\theta)$  there exists the logarithmic gradient (see [3]), and

$$\beta_{\theta}(x; z(x)) = \langle \lambda(\theta; x), z(x) \rangle + trz'(x).$$

This functional with respect to the continuity can be extended to smooth vector fields  $z(x): \aleph \to H$ .

**Example 2.** In the conditions of Example 1, we consider the vector field  $z(x): H_- \to H_-$ , possessing the bounded derivative:  $\sup_{x \in H_-} ||z'(x)|| < \infty$ .

If  $z: H_- \to H$ , then the logarithmic gradient exists. But if it is known additionally that  $z: H_- \to H_+$ , then  $\beta_{\theta}(x; z) = (\theta - x, z(x))_H + trz'(x)$ .

Here we cite some properties of the logarithmic derivative; their proof can be found in [2].

**Proposition 1.** Let the following conditions be fulfilled:

- (i) The measures  $P = P(\theta, \cdot)$  are differentiable along the vector  $h \in \aleph$ ;
- (ii) The functions f and g are differentiable along  $h \in \aleph$ ;
- (iii)  $f, g \in L_1(d_h P)$  and  $f'(x)h, g'(x)h \in L_1(P)$ ;
- (iv)  $(f'(x)h)g(x), f(x)(g'(x)h), f(x)g(x)\beta_{\theta}(x;h) \in L_1(P)$ . Then

$$\int_{\aleph} (f'(x)h)g(x)P(\theta;dx) = -\int_{\aleph} f(x)(g'(x)h)P(\theta;dx) = 
= -\int_{\aleph} f(x)g(x)\beta_{\theta}(x;h)P(\theta;dx).$$
(2)

**Proposition 2.** Let the measures  $P = P(\theta, \cdot)$  be differentiable along the vector  $h \in \aleph$  and the function  $\varphi(t) = \beta_{\theta}(x + th; H)$  be everywhere differentiable with  $\beta'_{\theta}(x; h)h \in L_2(P)$ . Then:

(i)  $P(\theta, \cdot)$  is twice differentiable along h;

(ii) 
$$d_{h,h}^2 P(\theta, \cdot) = \lfloor \beta'_{\theta}(x; h)h + (\beta_{\theta}(x; h))^2 \rfloor P(\theta; \cdot);$$
  
(iii)  $\int\limits_{\aleph} (\beta_{\theta}(x; h))^2 P(\theta; dx) = -\int\limits_{\aleph} \beta'_{\theta}(x; hP(\theta; dx).$ 

We will need measure smoothness with respect to the parameter, as well. Let, as above, we have the statistical structure  $\{\aleph, \Re(P(\theta; \cdot), \theta \in \Theta)\}$ , where  $\aleph$  is the separable, real Banach space, and let  $\Theta$  be a smooth manyfold imbedded into another separable, real Banach space  $\Xi$ . For any fixed  $A \in \Re$  and for the vector  $\vartheta \in \Xi$ , let us consider the derivative of the function  $\tau(\theta) = P(\theta; A)$  at the point  $\theta$  along  $\vartheta$ . We denote this derivative as follows:  $d_{\theta}P(\theta; A)\vartheta$ . For the fixed  $\theta$  and  $\vartheta$ , this derivative is the measure of alternating signs. It is easy to see that  $d_{\theta}P(\theta, \cdot)\vartheta << P(\theta, \cdot)$ , and by the Radon-Nikodym's theorem, there exists the measurable function  $l_{\theta}(x;\vartheta) = \frac{d_{\theta}P(\theta; dx)\vartheta}{P(\theta; dx)}$ .  $l_{\theta}(x;\vartheta)$  which is called logarithmic derivative of measure with respect to the parameter  $P(\theta; \cdot)$ .

When  $\Xi$  is the separable Hilbert space, by **K** we denote a space of measures for which the logarithmic derivative with respect to the parameter is representable in the form of a scalar product  $l_{\theta}(x; \vartheta) = (\mathbf{k}(x, \theta), \vartheta)_{\Xi}$ . In addition,  $\mathbf{k}(x, \theta)$  will be called a vector logarithmic gradient with respect to the parameter. For Examples 1 and 2,  $\lambda(x, \theta) = \theta - x$  and  $\mathbf{k}(x, \theta) = x - \theta$ .

For the family of measures  $(P(\theta;\cdot),\theta\in\Theta)$  possessing the logarithmic derivative with respect to the parameter along  $\theta$ , there exists the measure  $\nu$  dominating this family. It is known ([11]) that all measures  $P(\theta;\cdot)$  are mutually equivalent, and  $\frac{P(\theta_2;dx)}{P(\theta_1;dx)} = \exp\int_{\theta_1}^{\theta_2} l_{\theta}(x;\theta)d\theta$ .

#### 2. The Regularity Conditions

A statistical structure is a notion, derivative from the probabilistic space and from a random value. Therefore in some conditions of regularity the above two notions of logarithmic derivative should be connected. Here we point out these conditions (the conditions of regularity).

**Condition 1.**  $X(\theta) = X(\theta; \omega) : \Theta \times \Omega \to \aleph$ , and there exists the derivative  $X'(\theta)$  with respect to  $\theta$  along  $\theta \in \Xi_0$ , where  $\Xi_0 \subset \Xi$  is the subspace of  $\Xi$ . This derivative is a linear mapping  $\Xi \to \aleph$  for every  $\theta \in \Theta$ . For any  $\theta \in \Xi_0$  and  $\theta \in \Theta$ , we have  $||\Xi'(\theta)\theta||_{\aleph} \in L_2(\Omega, P)$ .

Condition 2.  $E\{X'(\theta)\vartheta|X(\theta)=x\}$  is strongly continuous as the function x for all  $\vartheta \in \Xi_0, \ \theta \in \Theta$ .

Condition 3. The family of measures  $(P(\theta;\cdot), \theta \in \Theta)$  possess the logarithmic derivative with respect to the parameter along constant directions from the dense in  $\Xi$  subspace  $\Xi_0 \subset \Xi$ , and  $l_{\theta}(x; \theta) \in L_2(\aleph, P(\theta)), \theta \in \Xi_0$ ,  $\theta \in \Theta$ .

**Condition 4.** The family of measures  $(P(\theta; \cdot), \theta \in \Theta)$  possess the logarithmic derivative along constant directions from the dense in  $\aleph$  subspace  $\aleph_0 \subset \aleph$  and  $\beta_{\theta}(x; h) \in L_2(\aleph, P(\theta)), h \in \aleph_0, \theta \in \Theta$ .

**Lemma 1.** Under the conditions of regularity 1)-4), for the logarithmic derivatives  $\beta_{\theta}(x;h)$  and  $l_{\theta}(x;\vartheta)$ , the equality

$$l_{\theta}(x; \vartheta) = -\beta_{\theta}(x; \mathbf{K}_{\theta, \vartheta}(x)) \text{ where } \mathbf{K}_{\theta, \vartheta}(x) = E\left\{\frac{d}{d\theta}X(\theta)\vartheta | X(\theta) = x\right\}$$
(3)

holds.

*Proof.* By the definition,  $P(\theta; A) = P(X^{-1}(\theta; A))$ . Let f(x) be the bounded, continuous differentiable along  $h \in \aleph$  real-valued function. By the change of variable formula,

$$\int_{\mathbb{R}} f(x)P(\theta; dx) = Ef(X(\theta)).$$

We differentiate both parts with respect to  $\theta$  along  $\vartheta$ . Thus we obtain

$$\int_{\mathbb{R}} f(x)d_{\theta}P(\theta;dx)\vartheta = E\frac{d}{dx}f(X(\theta))\frac{d}{d\theta}X(\theta)\vartheta,$$

or

$$\int\limits_{\aleph} f(x)l_{\theta}(x;\vartheta)P(\theta;dx) = \int\limits_{\aleph} f'(x)E\{X'(\theta)\vartheta|X(\theta) = x\}P(\theta;dx).$$

Denote 
$$K_{\theta,\vartheta}(x) = E\left\{\frac{d}{d\theta}X(\theta)\vartheta|X(\theta) = x\right\}$$
, and write

$$\int_{\mathbb{R}} f'(x)K_{\theta,\vartheta}(x)P(\theta;dx) = -\int_{\mathbb{R}} f(x)\beta_{\theta}(x,K_{\theta,\vartheta}(x))P(\theta;dx).$$

Since f(x) is arbitrary, we obtain (3).

**Example 3.** Let  $\Xi = R^2$ ,  $\Theta = R \times (0, \infty)$ ,  $\Xi_0 = R^2$ ,  $\aleph = R$ . For a random  $X(\theta_1, \theta_2)$ , the distribution  $P(\theta; A)$ ,  $A \in \mathbf{B} = \Re$ , where  $\mathbf{B}$ , the Borel  $\sigma$ -algebra in R, is prescribed by density

$$P(\theta_1, \theta_2; A) = \frac{1}{\sqrt{2\pi}\theta_2} \int_A \exp\left\{-\frac{(x - \theta_1)^2}{2\theta_2^2}\right\} dx.$$

Then for any  $\vartheta = (\vartheta_1, \vartheta_2)^T$ , we have

$$l_{\theta}(x:\vartheta) = \frac{(x-\theta_1)\theta_2\vartheta_1 - \theta_2^2\vartheta_2 + (x-\theta_1)^2\vartheta_2}{\theta_2^3}.$$

On the other hand, if by N we denote a standard normal distribution in R, then we can write  $X(\theta_1, \theta_2) = \theta_2 N + \theta_1$ . Hence

$$X'(\theta)\vartheta = \vartheta_1 + \frac{X(\theta) - \theta_1}{\theta_2}\vartheta_2.$$

Respectively,  $E\{X'(\theta)\vartheta|X(\theta)=x\}=\vartheta_1+\frac{x-\theta_1}{\theta_2}\vartheta_2=z(x)$ . In this situation,  $\lambda(\theta;x)=\frac{\theta_1-x}{\theta_2^2}$  and  $\mathbf{k}(x,\theta)=\left(\frac{x-\theta_1}{\theta_2^2},\frac{(x-\theta_1)^2-\theta_2^2}{\theta_2^3}\right)^T$ . Clearly, we have

$$\beta_{\theta}(x;z(x)) = \frac{\theta_1 - x}{\theta_2^2} \left( \vartheta_1 + \frac{x - \theta_1}{\theta_2} \vartheta_2 \right) + \frac{\vartheta_2}{\theta_2} = -l_{\theta}(x;\vartheta).$$

**Example 4.** Let  $X(\theta)$  be a normally distributed random element with values in the separable, real Hilbert space H with 0 mean and with kernel correlation operator  $\theta = B$ . Let  $P(\theta)$  be the corresponding Gaussian measure in H. In this case,  $\aleph = H$ ,  $\Re$ , is the Borel  $\sigma$ -algebra in it,  $\Xi = \mathbf{L}_1(H, H)$  is the Banach space of kernel operators in H with the norm  $||K||_1 = trK$ ,  $\Theta \subset \mathbf{L}_1(H, H)$  is the space of linear operators C such that  $B^{-1/2}CB^{-1/2}$  is the kernel operator in H. We have  $X(B) = B^{1/2}N$ , where N is the canonical Gaussian value in H. As a direction, we take the operator  $C \in \mathbf{L}_{-1}(H, H)$ . Calculations provide us with  $X'(B)C = 1/2 \left(B^{-1/2}CB^{-1/2}(X(B))\right)$  and  $E\{X'(B)C|X(B) = 1/2B^{-1/2}CB^{-1/2}x = z(x)$ . Respectively,

$$l_B(x;C) = -\beta_B(x;z(x)) = \frac{1}{2} \left( B^{-1/2}CB^{-1/2}x, B^{-1}x \right)_H - \frac{1}{2} tr_H B^{-1/2}CB^{-1/2}.$$

#### 3. The Cramer-Rao Inequality

Let  $\{\aleph, \Re, (P(\theta, \cdot), \theta \in \Theta)\}$  be the statistical structure corresponding to a random element  $X(\omega) = X(\theta, \omega)$ . Here,  $\aleph$  is the separable, real, reflexive Banach space,  $\Re$  is the  $\sigma$ -algebra of Borel sets, and  $\Theta \subset \Xi$  is an open subset of the separable, real Banach space  $\Xi$ . The conditions of regularity 1)–4) will be assumed to be fulfilled.

Let  $g(\theta) = E_{\theta}(T(X))$ , where  $T : \aleph \to R$  is a measurable mapping (statistics). For the statistics we take one more condition of regularity.

Condition 5. For the statistics  $T = T(x) : \aleph \to R$ , the equality

$$d_{\vartheta} \int_{\aleph} T(x)P(\theta; dx) = \int_{\aleph} T(x)d_{\vartheta}P(\theta; dx)$$

is valid.

**Theorem 1** (The Cramer-Rao Inequality). Let the conditions of regularity 1.–5. be fulfilled. Then

$$VarT(X) \ge \frac{(g'_{\vartheta}(\theta))^2}{E_{\theta}l_{\theta}^2(X;\vartheta)}.$$
 (4)

*Proof.* Having differentiated  $g(\theta)$  along  $\theta \in \Xi$ , we get

$$d_{\vartheta}E_{\theta}T(X) = d_{\vartheta} \int_{\aleph} T(x)P(\theta; dx) = \int_{\aleph} T(x)d_{\vartheta}P(\theta; dx)\vartheta =$$
$$= \int_{\aleph} T(x)l_{\theta}(x, \vartheta)P(\eta; dx) = E_{\theta}T(X)l_{\theta}(X; \vartheta).$$

Thus

$$d_{\vartheta}E_{\theta}T(X) = E_{\theta}T(X)l_{\theta}(X;\vartheta). \tag{5}$$

In (5), we put T(X) = 1. We obtain  $E_{\theta}l_{\theta}(X; \vartheta) = 0$ . Therefore

$$d_{\vartheta}E_{\theta}(T(X)) = E_{\theta}((T(X) - g(\theta))l_{\theta}(X; \vartheta)),$$

and hence

$$(d_{\vartheta}E_{\theta}(T(X)))^{2} \leq E_{\theta} (T(X) - g(\theta))^{2} \cdot E_{\theta}l_{\theta}^{2}(X; \vartheta)$$

which yields

$$VarT(X) \ge \frac{(g'_{\vartheta}(\theta))^2}{E_{\theta}l_{\theta}^2(X;\vartheta)}.$$

Corollary 1. Taking into account our Lemma, inequality (4) takes the form

$$VarT(X) \ge \frac{(g'_{\vartheta}(\theta))^2}{E_{\theta}\beta_{\theta}^2(X; E(X'(\theta)\vartheta|X))}.$$

**Example 5.** Let there be observed an  $X_k(\theta_1, \theta_2) = \frac{1}{\theta_1} e_k + \theta_2$  random value, where every  $e_k$  (k = 1, 2, ..., n) is an exponentially distributed random value with distribution density

$$p(x) = \begin{cases} e^{-x}, & x \ge 0\\ 0, & x < 0 \end{cases}.$$

Let  $X_1, X_2, \ldots, X_n$  be observations;  $X(\theta) = (X_1(\theta_1, \theta_2), \ldots, X_n(\theta_1, \theta_2))$ .

$$X'(\theta) = \begin{pmatrix} -\frac{1}{\theta_1^2} e_1 & 1\\ \dots & \dots & \dots\\ -\frac{1}{\theta_1^2} e_n & 1 \end{pmatrix} = \begin{pmatrix} -\frac{X_1(\theta) - \theta_2}{\theta_1} & 1\\ \dots & \dots & \dots\\ -\frac{X_n(\theta) - \theta_2}{\theta_1} & 1 \end{pmatrix}, \ \theta = (\theta_1, \theta_2).$$

If we choose the direction  $\vartheta = \begin{pmatrix} \vartheta_1 \\ \vartheta_2 \end{pmatrix}$ , then

$$X'(\theta)\vartheta = \begin{pmatrix} -\frac{X_1(\theta) - \theta_2}{\theta_1}\vartheta_1 + \vartheta_2 \\ \cdots \\ -\frac{X_n(\theta) - \theta_2}{\theta_1}\vartheta_1 + \vartheta_2 \end{pmatrix}.$$

Therefore

$$z(x) = E\{X'(\theta)\vartheta|X(\vartheta) = x\} = \begin{pmatrix} -\frac{x_1 - \theta_2}{\theta_1}\vartheta_1 + \vartheta_2\\ \cdots\\ -\frac{x_n - \theta_2}{\theta_1}\vartheta_1 + \vartheta_2 \end{pmatrix}.$$

For the exponential distribution X, the logarithmic derivative is

$$\lambda(x) = -\frac{1}{\theta_1} I \ (x \ge \theta_2). \text{ If } \Lambda = \begin{pmatrix} \lambda(x) \\ \cdots \\ \lambda(x) \end{pmatrix}, \text{ then }$$

$$\beta_{\theta}(x;h) = (\Lambda,h)_{R^n} = -\frac{I(x \ge \theta_2)}{\theta_1} \sum_{k=1}^n h_k.$$

Finally,

$$\beta_{\theta}(x; z(x)) = -\frac{I(\min x_k \ge \theta_2)}{\theta_1} \sum_{k=1}^n \left[ t - \frac{x_k - \theta_2}{\theta_1} \vartheta_1 + \vartheta_2 \right] +$$

$$+tr \begin{pmatrix} -\frac{\vartheta_1}{\theta_1} & 0 & \cdots & 0 \\ \cdots & \cdots & \cdots & \cdots \\ 0 & 0 & \cdots & -\frac{\vartheta_1}{\theta_1} \end{pmatrix} = \frac{I(\min x_k \ge \theta_2)}{\theta_1^2} \sum_{k=1}^n x_k -$$

$$-\frac{nI(\min x_k \ge \theta_2)\theta_2\vartheta_1}{\theta_1^2} - \frac{nI(\min x_k \ge \theta_2)\vartheta_2}{\theta_1} - \frac{n\vartheta_1}{\theta_1}.$$

## 4. The Method of Maximal Likelihood

Let  $\{\aleph, \Re, (P(\theta), \theta \in \Theta)\}$  be the statistical structure corresponding to a random element  $X = X(\theta) = X(\theta, \omega), \ \omega \in \Omega$ , where  $\aleph$  is the separable, real Banach space,  $\Re$  is the  $\sigma$ -algebra of the Borel subsets, and  $\Theta$  is an open subset of another separable, real Banach space  $\Xi$ . We assume that the family of measures  $(P(\theta,\cdot),\theta \in \Theta)$  possess the logarithmic derivative with respect to the parameter  $l_{\theta}(x,\vartheta)$  along  $\vartheta \in \Xi$ . Then by Theorem 1, there exists the logarithmic derivative with respect to the measure, and  $l_{\theta}(x,\vartheta) = -\beta_{\theta}(x,E\{X'(\theta)\vartheta|X(\theta) = x\}$ .

Consider a structure of repeated sampling  $\{\aleph, \Re, (P(\theta), \theta \in \Theta)\}^n = \{\aleph^n, \Re^n, (P^n(\theta), \theta \in \Theta)\}.$ 

**Theorem 2.** If there exists the logarithmic derivative  $l_{\theta}(x, \vartheta)$  with respect to the parameter in the statistical structure  $\{\aleph, \Re, (P(\theta), \theta \in \Theta)\}$ , then there likewise exists the logarithmic derivative  $L_{\theta}((x_1, \ldots, x_n), \vartheta^n)$  with respect to the parameter, along  $\vartheta^n \stackrel{def}{=} (\vartheta, \ldots, \vartheta)$ , for the structure of repeated sampling  $\{\aleph, \Re, (P(\theta), \theta \in \Theta)\}^n$ , and we have

$$L_{\theta}((x_1, \dots, x_n), \vartheta^n) = \sum_{k=1}^n l_{\theta}(x_k, \vartheta) =$$

$$= -\sum_{k=1}^n \beta_{\theta} (x_n, E\{X_k'(\theta)\vartheta | X_k(\theta) = x_k\}). \tag{6}$$

*Proof.* Since by the condition, there exists  $d_{\theta}^{\vartheta}P(\theta)$ , it is not difficult to calculate that there likewise exists  $d_{\theta}^{\vartheta,\dots,\vartheta}P^n(\theta)=\sum_{k=1}^n d_{\theta}^{\vartheta}(x_k,\vartheta)\prod_{\substack{j=1\\j\neq k}}^n P(\theta)$ . This

last one is absolutely continuous with respect to  $P^n(\theta)$ . By the Radon-Nikodym's theorem, we find that the statements of the theorem, as well as formula (6), are valid.

It follows from the theorem that in the case under consideration we can formulate the principle of maximal likelihood.

Let  $X_1, X_2, ..., X_n$  be sampling from the random value  $X(\theta)$ , where  $\theta$  is an unknown parameter to be estimated by means of sampling. Assume that there exists the logarithmic derivative  $l_{\theta}(x, \theta)$  with respect to the parameter, along any vector  $\theta \in \Xi_0$ , of distribution  $P(\theta)$ , which corresponds to  $X(\theta)$  and has the form  $l_{\theta}(x, \theta) = \langle \lambda(x, \theta), \theta \rangle$ . Here,  $\Xi_0$  is the dense subset of  $\Xi$ .

As is known, all measures  $P(\theta)$  are equivalent to each other. Let  $\theta_0 \in \Xi_0$  be a fixed point. Consider the likelihood function

$$\frac{dP(\theta)}{dP(\theta_0)}(x) = \rho(x,\theta).$$

It can be easily seen that for  $P \in \mathcal{L}$  the above equality results in

$$\frac{\rho_{\theta}'(x,\theta)\vartheta}{\rho(x,\theta)} = l_{\theta}(x,\vartheta).$$

For the sampling  $X_1, X_2, \dots, X_n$ , the likelihood function is

$$L(X_1, X_2, \dots, X_n, \theta; \vartheta) = \prod_{k=1}^n \rho(X_k, \theta).$$

According to the likelihood principle, the estimate of maximal likelihood will be called the value  $\theta = \hat{\theta}$  which supplies the likelihood function L with

maximum (provided that such value of the parameter  $\theta$  exists). Since

$$\ln L(X_1, X_2, \dots, X_n, \theta; \vartheta) = \sum_{k=1}^n \ln \rho(X_k, \theta),$$

the condition for maximum allows us to formulate this definition in terms of the logarithmic derivative with respect to the parameter:

The estimate of maximal likelihood  $\widehat{\theta}$  with respect to the direction  $\vartheta$  is called the root (if exists) of the equation

$$\sum_{k=1}^{n} l_{\theta}(x_k, \vartheta) = 0, \quad \forall \vartheta \in \Xi$$
 (7)

with respect to  $\theta$ , under the condition that the expression  $\frac{d}{d\theta}l(x,\theta)$  is defined negatively.

By Lemma 1, equation (7) can be replaced by

$$\sum_{k=1}^{n} \{ \mathbf{k}(x_k, \theta), K(x_k, \theta)\vartheta \}_{\aleph} + trK_x'(x_k, \theta)\vartheta \} = 0, \quad \forall \vartheta \in \Xi_0.$$
 (8)

 $x_k$  in formulas (7) and (8) are the values of Xi, experimentally.

**Example 6.** Let in the equipped Hilbert space  $H_+ \subset H \subset H_-$  be considered the sampling  $X_1, X_2, \ldots, X_n$  from the canonical Gaussian value with an unknown mean  $\theta$ , for which  $\beta_{\theta}(x,h) = (\theta-x,h)_H$ ,  $h \in H_+$ . Clearly,  $X(\theta) = N + \theta$ , where N is the canonical Gaussian value with zero mean.  $X'(\theta) = I$ ,  $X'(\theta)h = h$ , and hence  $E\{X'_k(\theta)h|X_k(\theta) = x\} = h$ . Thus, (8) takes the form

$$\sum_{k=1}^{n} (\theta - x_k, h)_H = 0,$$

whence

$$(\widehat{\theta}, h)_H = \frac{1}{n} \sum_{k=1}^n (X_k, h)_H \text{ and } \widehat{\theta} = \frac{1}{n} \sum_{k=1}^n X_k = \overline{X}^1.$$

In addition, 
$$\frac{1}{n} \sum_{k=1}^{n} \frac{d^h}{d\theta} (x - \theta, h)_H = -n||h||_H^2 \le 0.$$

$$(j^*\widehat{\theta},jh)_H = \frac{1}{n}\sum_{k=1}^n (j^*X_k,jh)_H, \ \text{ and again } \ \widehat{\theta} = \frac{1}{n}\sum_{k=1}^n X_k = \overline{X}.$$

If  $h \in H$ , then, as is known, the expression  $\beta_{\theta}(x,h) = (\theta - x,h)_H$  extends with respect to the continuity as a measurable linear functional. The obtained in such a way so-called Skorokhod's integral solves our problem.

<sup>&</sup>lt;sup>1</sup> This equality is obtained under the condition  $h \in H_+$  and  $\theta - x_k \in H_-$ . If j is the embedding operator  $j: H_+ \to H$  and  $j^*: H_- \to H$ , we can rewrite

As an application, we consider a random process  $x(t) = \varphi(t) + w(t)$ , where w(t) is a standard Wiener's process,  $\varphi \in C[,\infty) = \Xi$  is an unknown component of the observable process. Surely,  $Ex(t) = \varphi(t)$ . In this case,  $H_+ = C'[0,\infty)$ ,  $H_- = L_2[0,\infty)$ .

If 
$$x_1(t), x_2(t), \dots, x_n(t)$$
 are observed, then  $\widehat{\varphi}(t) = \frac{1}{n} \sum_{k=1}^{n} x_k(t)$ .

If  $H = \mathbb{R}^n$  is finite-dimensional, we obtain the estimate of maximal likelihood along any vector  $h = (h_1, \dots, h_n)$ :

$$\widehat{\theta} = (\widehat{\theta}_1, \dots, \widehat{\theta}) = \left(\frac{1}{n} \sum_{k=1}^n X_k^1, \dots, \frac{1}{n} \sum_{k=1}^n X_k^m\right).$$

#### 5. Consistensity of the Maximal Likelihood Estimate

Let the statistical structure  $\{\aleph, \Re, (P(\theta), \theta \in \Theta)\}$  possess the logarithmic derivative with respect to the parameter along the constant vector  $\vartheta \in \Xi_0 - l_\theta(x, \vartheta)$ . We introduce a Kulbak-Leybler type function of distance for pairs of measures:

$$D(\theta_1, \theta_2) = E_{\theta_1} \{ l_{\theta_1}(x, \theta_2 - \theta_1) - l_{\theta_2}(x, \theta_2 - \theta_1) \}.$$
 (9)

**Example 7.** In the equipped Hilbert space  $H_+ \subset H \subset H_-$ , for canonical Gaussian measures  $\mu_1$  and  $\mu_2$  with, respectively,  $\theta_1$  and  $\theta_2$  means, the distance is  $D(\mu_1, \mu_2) = (\theta_1 - \theta_2, \theta_2 - \theta_1)_H = -||\theta_2 - \theta_1||_H^2$ .

**Lemma 2.** Let the family be defined uniquely by the parameter, i.e., if  $P(\theta_1) = P(\theta_2)$ , then  $\theta_1 = \theta_2$ , and vice versa. And if  $D(\theta_1, \theta_2) \geq 0$ , then  $P(\theta_1) = P(\theta_2)$ , and vice versa.

*Proof.* We note immediately that  $D(\theta_1, \theta_2) \leq 0$ . Indeed,

$$\begin{split} D(\theta_1, \theta_2) &= E_{\theta_1} \{ l_{\theta_1}(x, \theta_2 - \theta_1) - l_{\theta_2}(x, \theta_2 - \theta_1) \} = \\ &= E_{\theta_1} (\lambda(X, \theta_1) - \lambda(X, \theta_2), \theta_2 - \theta_1)_{\Xi} = \\ &= E_{\theta_1} (\lambda'_{\theta}(X, \theta_1 + \tau(\theta_2 - \theta_1))(\theta_2 - \theta_1), \theta_2 - \theta_1)_{\Xi} \le 0, \ (0 \le \tau \le 1). \end{split}$$

Thus we can see that if  $D(\theta_1, \theta_2) \ge 0$ , then  $\theta_1 = \theta_2$ , which implies  $P(\theta_1) = P(\theta_2)$ , and vice versa.

**Theorem 3.** If  $\Xi_0$  is a convex precompact set, then the estimate of maximal likelihood is consistent.

*Proof.* Let  $\widehat{\theta}$  be a solution of equation (8) or (9). Consider the difference

$$\varphi_n(t) = \frac{1}{n} \ln L(x_1, \dots, x_n, \theta; t(\widehat{\theta} - \theta)) - \frac{1}{n} \ln L(x_1, \dots, x_n, \widehat{\theta}; t(\widehat{\theta} - \theta)).$$

Since  $\hat{\theta}$  for any t is a point of maximum,  $\varphi_n(t)$  decreases with respect to t on [0,1]. On the other hand, by the strong law of large numbers,

$$\varphi_n(t) = \frac{1}{n} \ln \frac{L(x_1, \dots, x_n, \theta; t(\widehat{\theta} - \theta))_{a.s}}{L(x_1, \dots, x_n, \widehat{\theta}; t(\widehat{\theta} - \theta))} \to \varphi(t),$$

and  $\varphi_n(t)$  is likewise the decreasing function. Therefore,  $\varphi'(t) \leq 0$ . But

$$\varphi'(t) = P \lim_{n} \varphi'_{n}(t) \text{ and } \varphi'_{n} \stackrel{a.s}{\to} E_{\widehat{\theta}}\{l_{\theta}(X, \widehat{\theta} - \theta) - l_{\widehat{\theta}}(X, \widehat{\theta} - \theta)\} = D(\widehat{\theta}, \theta) \leq 0.$$

By Lemma 2, we have 
$$\theta = \widehat{\theta}$$
.

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