

Ergodic Theory in Classical and Bayesian inference

Artur Oscar Lopes

Inst. Mat. Est. - UFRGS

Av. Bento Goncalves 9500 - Porto Alegre - 91500-000 - Brazil

June 5, 2026

Abstract

We begin by presenting the mathematical rationale underlying classical deductive inference. We then introduce the foundational ideas of the Bayesian inference framework. Results lying at the interface of Statistics and Ergodic Theory are outlined, providing a theoretical framework applicable to the prediction and analysis of real-world phenomena from random data. This text is expository in nature — no new results are presented; rather, recent published results are described in a didactic manner. Throughout, we work with Hlder equilibrium measures, which encompass a substantially more general class of processes than i.i.d. ones.

Keywords: Classical Inference, Bayesian inference, Hölder equilibrium probability, Kullback-Leibler divergence, loglikelihood ratio, loss function, prior and posterior probability

Mathematics Subject Classification: 37D35; 62F05; 62F15

email: arturoscar.lopes@gmail.com

1 Inductive Inference

This short note is intended to explain to mathematicians some applications of Ergodic Theory to some specific problems in Classical and Bayesian inference. We believe there is a need for an exposition that deals with inference in a didactic way and without major technicalities. In principle, the Bayesian framework is opposed to the frequentist point of view. More explanations are needed on this issue.

In Section 2 we will provide some simple examples for the reader that is not familiar with the topic of inference.

At the end of the paper, we will consider inference results for Hölder equilibrium probabilities (see [28]), which is a class much more general than the classical i.i.d. processes; any shift invariant probability can be weakly approximated by a Hölder equilibrium probability (see [18] or [14]).

Recently the authors A. Nobel, K. McGoff, S. Mukherjee, and N.S. Pillai presented several interesting results in the interplay of Statistics and Ergodic Theory. We will describe in a didactic and synthetic way another line of investigation: we highlight and put in context some new results by M. Denker, H. H. Ferreira, M. Kesseböhmer, A. O. Lopes, S. R. C. Lopes, J. Mengue, and P. Varandas

In the frequentist setting, one assumes that the probability of an event is proportional to the times that the event occurs in a sequence of observations of random data; in this way, the Ergodic point of view is quite pertinent. The section 2 will describe basic results for classical inference under such a point of view.

Frequentists say the data are random and the parameters are fixed; on the other hand, Bayesians say the data are fixed, and the parameters are random.

In real-world problems, sometimes a satisfactory sequence of data is available for a frequentist analysis, and in some cases, not. The choice of a purely Bayesian, or purely frequentist analysis depends very much on the type of the problem, and the conviction of the investigator; it also depends on the relevance or not of the available data (a subjective matter) for the formulation of the Statistical model in question.

In the frequentist framework, typically, there is a fixed probability that generates a time series of random data. On the other hand, in the Bayesian setting, one can consider randomness in the probabilities that are relevant to the problem under consideration, or that the time series that eventually exists is small, or nonexistent, and thus of little relevance. In several cases, some subjectivity in the formulation of the model is necessary.

From section 6.1 in the book [1] we get the claim:

“Although the two approaches are conceptually different, they are nevertheless not *complete strangers* to each other and may benefit from cooperation. Frequentist analysis of Bayesian procedures can be very useful for better understanding and validating their results.”

In the present work, we consider some aspects of the Bayesian framework where one can take advantage of Ergodic Theory. The final results will be produced by mixing these two apparently opposed points of view. This will

be considered in Section 3.

What is Inference in general terms?

Quoting Chapter 1 in [3]

The process of drawing conclusions from available information is called inference. When the available information is sufficient to make unequivocal, unique assessments of truth, we speak of making deductions: based on a certain piece of information, we deduce that a certain proposition is true. The method of reasoning leading to deductive inferences is called logic. Situations where the available information is insufficient to reach such certainty lie outside the realm of logic. In these cases, we speak of doing inductive inference, and the methods deployed are those of probability theory and entropic inference.

Given a certain problem to be analyzed, the use of the frequentist or Bayesian point of view may be a matter of dispute among different researchers. The use of each of the two options and their consequent predictions in comparison with the real world may show in each case whether it is more relevant to use one or the other.

Note that in the formulation of a statistical model, even when there are preliminary historical data, their quantity can be understood as sufficient (a subjective question) for the corresponding analysis or not. In the latter case, a Bayesian researcher often uses simulations with random data, obtained via a computer, to calibrate, or validate, the model, and/or, to prove its efficiency (see [15]). In this case, it can be said that there is a certain frequentist remnant. In some cases, such as when there is a long dependency in the data, obtaining long time series via random simulations on a computer may be time-consuming (see, for instance, [24], [25], [30], and [29]).

We do not intend to exhaust all the possible topics covered by Bayesian statistics, but only to describe in an intelligible mathematical way some aspects that may allow the reader to understand some of its essential ingredients (see beginning of Section 3). For most of the models in Bayesian Statistics, the Ergodic point of view is not so relevant (see for instance, [32], [1], and [2]).

Some results relating Ergodic Theory and Statistics are [5], [7], [8], [9], [10], [11], [12], [13], [4], [17], [19], [20], [21], [26], [27], [16], [22], [23] and [6].

In [3] the author explores the point of view of deriving the main postulates of Physics from inference.

2 Classical Inference

An appropriate form of describing states of partial belief (or uncertainty) is via probabilities q on a certain set X . A quite natural strategy is the following: the state of partial belief q^0 should be updated to a new state of partial belief p^0 when new information is available. The mathematical procedure to get the posterior probability p^0 is based on a mathematical formalism that takes into account the prior probability q^0 and also a function l , which is usually called the loss function (sometimes called risk function). The loss function may change according to the specific problem one is considering. We say that p^0 was obtained via an inference process obtained from q^0 and l .

The new information may arrive in the form of a random time series, but not necessarily in this form. In some cases, this information has to be processed via a minimization procedure involving probabilities.

The states of partial belief may be indexed by a set Θ of possible probabilities. This way, we get a family of probabilities $p(\theta)$ for $\theta \in \Theta$. The probabilities q^0 and p^0 are among this set of possible probabilities.

The randomness of the probability q^0 may be caused by lack of information, and inductive inference is commonly used in Information theory.

A simple example of inference can be described in the following way. Consider the case where $X = \{1, 2, \dots, d\}$ and

$$\Theta = \{p = (p_1, p_2, \dots, p_d) \mid \sum_{j=1}^d p_j = 1, p_j \geq 0, j = 1, 2, \dots, d\}.$$

Probabilities on Θ are denote by $p(\theta)$. The parameter θ is indexed by (p_1, p_2, \dots, p_d) .

We denote by Ω the symbolic space $X^{\mathbb{N}}$. Each probability $p(\theta)$ is associated with i.i.d. Stochastic Process, which means a probability μ_θ on Ω .

We consider a **special point** $\theta_0 \in \Theta$ **which will be responsible for producing random sequences** $y = (y_1, y_2, \dots, y_n, \dots) \in \Omega$. The associated probability on Ω is denoted by μ_{θ_0} , which in principle we don't know. However, we can observe finite random samples over time. Given a sample sequence $y = (y_1, y_2, \dots, y_n, \dots) \in \Omega$, we denote y^n the finite word $y^n = (y_1, y_2, \dots, y_n) \in X^n$. Moreover, we denote by $C_n(y)$ the cylinder $\overline{y_1, y_2, \dots, y_n} = \overline{y^n}$.

A very useful transformation in the study of probability measures in the symbolic space Ω is the shift transformation $\sigma : \Omega \rightarrow \Omega$. We set

$$\sigma(y_1, y_2, \dots, y_n, \dots) = (y_2, y_3, \dots, y_n, \dots).$$

An Stochastic Process with values on X determine a probability P on Ω . A stochastic process being stationary corresponds to the property that for any cylinder $C_n(y)$, it holds that $P(\sigma^{-1}(C_n(y))) = P(C_n(y))$. If this property is true we say that P is invariant for the action of the shift.

In most cases, for applications in the real world, a fixed random source is the object of interest in the problem under consideration, which is not precisely known but generates random samples.

The finite sample $y^n = (y_1, y_2, \dots, y_n)$ corresponds to the new information at time n .

We don't know which θ_0 produced the sample y , but if we have to bet, it would be better to choose the $\hat{\theta} = \hat{\theta}(y)$ such that

$$\mu_{\hat{\theta}}(C_n(y)) \geq \mu_{\theta}(C_n(y)),$$

when compared to other possible θ . In this case, it is natural to say that $\mu_{\hat{\theta}}$ is more suited to the data y^n than μ_{θ} (see Section 1.2 in [1]). We will elaborate on this claim.

The likelihood ratio

$$\frac{\mu_{\hat{\theta}}(C_n(y))}{\mu_{\theta_0}(C_n(y))}$$

shows the strength of the evidence in favor of $\hat{\theta} = \hat{\theta}(y)$ against θ_0 .

In this way is natural to try to maximize

$$\log \mu_{\hat{\theta}}(C_n(y)) - \log \mu_{\theta_0}(C_n(y)).$$

Then, it is natural to consider a minimization problem associated with the *loss functions* $\ell_n : \Theta \times \Omega \rightarrow \mathbb{R}$, $n \in \mathbb{N}$, given by

$$\ell_n(\theta, y) = -\log \left(\frac{\mu_{\theta}(C_n(y))}{\mu_{\theta_0}(C_n(y))} \right) = \log \left(\frac{\mu_{\theta_0}(C_n(y))}{\mu_{\theta}(C_n(y))} \right) \quad (1)$$

This loss function is of the log-likelihood form. We refer the reader to Chapter 2, Section 16 in [2] for general results on the maximum likelihood method.

At first, one can ask if there exists a conceptual issue with the above reasoning because we claim we do not know θ_0 , but ℓ_n encapsulates the information of θ_0 . We will show that this will not be a problem (see (9)).

For a given $n \in \mathbb{N}$, and y chosen at random with respect to θ_0 , we are interested in the value $\theta^n = \theta_y^n \in \Theta$ minimizing $\theta \rightarrow \ell_n(\theta, y)$. In this case $p(\theta^n)$ was obtained by inductive inference from l and $p(\theta_0)$.

We would like to show that for y chosen at random with respect to θ_0 , we get that $\theta^n \rightarrow \theta_0$ when $n \rightarrow \infty$. The bottom line is: from the loss function

l , by getting more and more information, we will be able to determine the source responsible for randomness.

Example 1. Set $\theta_0 = \tilde{p} = (1/3, 1/3, 1/3)$, and for fixed $n > 0$, consider a random sample $y = (y_1, y_2, \dots, y_n)$ of size n , obtained from the independent probability P on $\{1, 2, 3\}^{\mathbb{N}}$ associated to \tilde{p} . The observer knows in advance that the source is produced by an independent probability P but does not know the exact $\theta = (q_1, q_2, q_3)$, $q_1 + q_2 + q_3 = 1$, that produces it. From the information of the sample y he wants to obtain the best approximate values $\hat{\theta} = (\hat{q}_1, \hat{q}_2, \hat{q}_3)$. The probability \tilde{p} is unknown to the observer. In this example we follow the frequentist point of view.

As we obtain a sample of larger size n , one can obtain a better approximation of the true values of the unknown $\theta = (q_1, q_2, q_3)$ that determine the probability P . Given a sample y , what is the best guess? We believe this simple example will illustrate the effectiveness of the maximum likelihood estimator procedure for obtaining a good guess $\hat{\theta} = \hat{\theta}_n = (\hat{q}_1, \hat{q}_2, \hat{q}_3)$. We denote by $C_n(y)$ the cylinder set $\overline{y_1, y_2, \dots, y_n} \in \{1, 2, 3\}^{\mathbb{N}}$. The function g_n will be defined as the maximum likelihood estimator

$$g_n(p) = \frac{1}{n} \log \left(\frac{p(C_n(y))}{\tilde{p}(C_n(y))} \right).$$

For fixed n , one is interested in the probability $p_n = (p_n^1, p_n^2, p_n^3)$ maximizing g_n . The loss function is $-g_n$. Assume the string y has n_j occurrences of each symbol $j = 1, 2, 3$. In this case $n_1 + n_2 + n_3 = n$. Using Lagrange multipliers one can show that the probability p_n maximizing the above expression is such that

$$p_n^j = \frac{n_j}{n}, \quad j = 1, 2, 3. \quad (2)$$

We elaborate on this: there are two constraints $p_1 + p_2 + p_3 = 1$ and $n_1 + n_2 + n_3 = n$. The gradient of $G(p_1, p_2, p_3) = p_1 + p_2 + p_3$ is $(1, 1, 1)$. The gradient of

$$F(p_1, p_2, p_3) = \log(p_1)^{n_1} (p_2)^{n_2} (p_3)^{n_3} - \log 3^{-n} \quad (3)$$

is

$$\nabla F(p_1, p_2, p_3) = \left(\frac{n_1}{p_1}, \frac{n_2}{p_2}, \frac{n_3}{p_3} \right), \text{ and we assume it is equal to } (\lambda, \lambda, \lambda).$$

That is, the Lagrange multipliers method produce

$$\lambda = \frac{n_1}{p_1}, \quad \lambda = \frac{n_2}{p_2}, \quad \lambda = \frac{n_3}{p_3}. \quad (4)$$

From the above we derive

$$\lambda = \lambda(p_1 + p_2 + p_3) = n_1 + n_2 + n_3 = n. \quad (5)$$

Therefore, from (4) and (5) it follows that (2) holds.

One can show that $p_n \rightarrow \tilde{p}$ as $n \rightarrow \infty$. Indeed, for fixed j , denote $I_{\bar{j}} : \{1, 2, 3\}^{\mathbb{N}} \rightarrow \mathbb{R}$ the indicator function of the cylinder $\bar{j} \subset \{1, 2, 3\}^{\mathbb{N}}$. Note that, by the Law of Large Numbers (or by the Birkhoff Ergodic Theorem),

$$\frac{n_j}{n} = \frac{\sum_{k=0}^{n-1} I_{\bar{j}}(\sigma^k(y))}{n} \rightarrow \int I_{\bar{j}} dP = P(\bar{j}) = 1/3.$$

The sample is finite: suppose $y = (2, 1, 2, 1, 3, 2, 1, 3, 2, 3)$, where $n = 10$. The guess is

$$\hat{\theta} = (\hat{q}_1, \hat{q}_2, \hat{q}_3) = (3/10, 4/10, 3/10).$$

Note that in (3) the expression after the minus sign does not depend on p_1, p_2, p_3 . In this way, the reasoning above would work for a further problem with other values $\theta_0 = \tilde{p} \neq (1/3, 1/3, 1/3)$, producing different samples and

deriving others $\hat{\theta}$.

◇

Following [5], given the ergodic probability measures μ_{θ_0} and μ_{θ} , the relative entropy $h(\mu_{\theta} | \mu_{\theta_0})$ is given by the limit

$$\begin{aligned} h(\mu_{\theta_0} | \mu_{\theta}) &:= \lim_{n \rightarrow \infty} \frac{1}{n} \log \left(\frac{\mu_{\theta_0}(C_n(y))}{\mu_{\theta}(C_n(y))} \right) = \\ &\sum_{r=1}^d p(\theta_0)_r \log p(\theta_0)_r - \sum_{r=1}^d p(\theta)_r \log p(\theta_0)_r = \\ &\int p(\theta_0) \log p(\theta_0) - \int p(\theta) \log p(\theta_0) \geq 0, \end{aligned} \tag{6}$$

which exists and is non-negative, for μ_{θ_0} -almost every $y = (y_1, y_2, \dots, y_n, \dots) \in \Omega$. The minimum value of (6) is zero, and it is attained just when $\theta = \theta_0$.

Note that $h(\mu_{\theta_0} | \mu_{\theta}) = 0$, if and only if, $\theta = \theta_0$. In the present case, if $h(\mu_{\theta_0} | \mu_{\theta_n}) \rightarrow 0$, then $\theta_n \rightarrow \theta_0$, when $n \rightarrow \infty$.

Given θ_1, θ_2 , if

$$\int p(\theta_0) \log p(\theta_0) - \int p(\theta_1) \log p(\theta_0) > \int p(\theta_0) \log p(\theta_0) - \int p(\theta_2) \log p(\theta_0), \quad (7)$$

then μ_{θ_2} is preferred, when compared with μ_{θ_1} .

Note that (7) is equivalent to

$$-\int p(\theta_1) \log p(\theta_0) > -\int p(\theta_2) \log p(\theta_0), \quad (8)$$

and then, given the finite random sample $y^n = (y_1, y_2, \dots, y_n)$, it is natural to say, taking into account the limit, that for large n , μ_{θ_2} is preferred when compared with μ_{θ_1} , if

$$\mu_{\theta_2}(C_n(y)) \geq \mu_{\theta_1}(C_n(y)).$$

Then, given the random data $y^n = (y_1, y_2, \dots, y_n)$, from the inference point of view, we should look for the θ maximizing

$$\mu_{\theta}(C_n(y)) = \mu_{\theta}(\overline{y_1, y_2, \dots, y_n}). \quad (9)$$

A natural question is: in the long run, taking large n in the above procedure, in the limit, will we localize θ_0 ? We will see that the answer is yes.

In the maximum-likelihood estimator procedure (see definition 1 in Chapter II.16 on page 69 in [2]), for fixed θ_0 , we want to maximize in θ the value

$$\int p(\theta) \log p(\theta_0) - \int p(\theta_0) \log p(\theta_0) \leq 0, \quad (10)$$

which is equivalent to minimize in θ the expression (6).

The main question is: in the long run, ingetting larger and larger samples, can we identify μ_{θ_0} ?

Suppose $E \subset \Theta$ is a closed set such that $\theta_0 \notin E$, then, there exists an $\alpha > 0$, such that,

$$\inf \left\{ \sum_{r=1}^d p(\theta)_j \log p(\theta_0)_j - \sum_{r=1}^d p(\theta_0)_j \log p(\theta_0)_j \mid \theta \in E \right\} < -\alpha < 0. \quad (11)$$

We can choose as an example E of the form $\{\theta \mid |\theta - \theta_0| > \epsilon\}$, for small ϵ .

Then, for $\theta \in E$, and for μ_{θ_0} -almost every $y = (y_1, y_2, \dots, y_n, \dots) \in \Omega$, we get from (6), that for large n

$$-\log\left(\frac{\mu_\theta(C_n(y))}{\mu_{\theta_0}(C_n(y))}\right) \geq \alpha n. \quad (12)$$

Suppose F is a closed set of the form

$$\left\{\theta \mid -\frac{\alpha}{2} \leq \int p(\theta) \log p(\theta_0) - \int p(\theta_0) \log p(\theta) \leq -\frac{\alpha}{4}\right\}$$

not containing θ_0 . Then, for $\theta \in F$, and for μ_{θ_0} -almost every $y = (y_1, y_2, \dots, y_n, \dots) \in \Omega$, we get for large n

$$\frac{\alpha}{2} n \geq -\log\left(\frac{\mu_\theta(C_n(y))}{\mu_{\theta_0}(C_n(y))}\right) \geq \frac{\alpha}{4} n. \quad (13)$$

Note that $F \cap E = \emptyset$. Therefore, given n , for μ_{θ_0} -almost every $y = (y_1, y_2, \dots, y_n, \dots) \in \Omega$, when looking for a θ minimizing

$$\ell_n(\theta, y) = -\log\left(\frac{\mu_\theta(C_n(y))}{\mu_{\theta_0}(C_n(y))}\right),$$

the optimal θ will not be in the set E for sure. Indeed, taking large samples y^n , in the long range, a θ_2 in F is preferred when compared with a θ_1 in E (see (7)); this is so because, $\alpha n > \frac{\alpha}{2} n$.

In this way, the minimizing solution $\theta^n = \theta_y^n$, for the function $\theta \rightarrow \ell_n(\theta, y)$, for large n , will not be in the set E .

A full proof of the results described above follows from the reasoning of [20].

The relative entropy is also known as Kullback-Leibler divergence (see [13] and [21]).

3 Bayesian Inference

We present here a particular point of view concerning a certain class of problems in the topic, which is very broad; other frameworks are also relevant and interesting, but not covered here.

In the same way as in the last section, the states of partial belief are indexed by a set $\Theta = [0, 1]^n$. Now the elements in Θ are probabilities μ on $\Omega = \{1, 2, \dots, d\}^{\mathbb{N}}$. In this way, we get a family of probabilities μ_θ , for $\theta \in \Theta$. But now, in the Bayesian framework, we assume that the parameters

θ are in a state of partial belief which is described by a probability Q_0 in Θ . This probability Q_0 may be associated with complete a priori ignorance, or alternatively, describe some subjective belief of those who wish to apply the method.

I. In one possible setting, there exists a special unknown parameter θ_0 , such that μ_{θ_0} produces the relevant samples, in a similar way as in the last section; but, now we set an initial distribution of probability given by a certain Q_0 on Θ , describing the uncertainty on θ_0 .

II. In another possible setting, each $\mu_{\theta}, \theta \in \Theta$, produces samples. In this case, the role of the probability Q_0 is to average the randomness of the samples obtained from the random family $\mu_{\theta}, \theta \in \Theta$.

In any case, in principle, Q_0 describes fuzzy information on probabilities. From different types of inference procedures, given a loss function l , not necessarily of log-likelihood type, one can get a posterior probability Q_1 in Θ , from the prior probability Q_0 .

The loss function l can have different forms depending on the problem of interest, and l may incorporate, or not, some subjective belief of those who wish to apply the method. This applies to both problems: types I and II.

One should consider two more alternatives that apply to cases I and II: a) whoever wishes to apply the method has a pertinent finite random sample of arbitrarily long size at their disposal; b) the item a) is not fulfilled.

In the first case a), the frequentist point of view could be incorporated into the Bayesian framework. Then, taking into account the sequence of finite random data $y^n = (y_1, y_2, \dots, y_n)$, $n \in \mathbb{N}$, and a family of loss functions l_n , we can obtain a family of probabilities Q_n in Θ . In this case, the information of Q_n , $n \in \mathbb{N}$, is less uncertain than the one given by Q_0 . We are interested in the eventual limit of such a family Q_n on Θ , when $n \rightarrow \infty$.

In case b), using an specific l , one gets the posterior Q_1 from the prior Q_0 , and this should satisfy (by particular reasons inherent to the problem) the properties aimed at by the one who applies the method. In this case, in the establishment of the Statistical model, a greater degree of subjectivity is inherent in its formulation. The choice of l needs to incorporate the appropriate relevance that is natural to the problem under consideration

Here we choose to describe the case where we assume I) and a). Subjectivity in the formulation of this model is not necessary. Other cases are of interest, of course, and we refer the reader to [1] and [2].

We point out that the Bayes formalism, and the Bayes rule, can be deduced via an Optimal Information procedure, as described in [31]; the topic was the subject of a more in-depth study, and considered in a more general setting in [19].

Assume that the probability μ_θ on $\Omega = \{1, 2, \dots, d\}^{\mathbb{N}}$, is the equilibrium probability for a Hölder potential $A_\theta : \Omega \rightarrow \mathbb{R}$ (see [28]) indexed by a parameter $\theta \in \Theta = [0, 1]^k$ (see [27], [11] and [20]). We will mention some new inference results concerning such a family.

For a special fixed parameter $\theta_0 \in \Theta$ we will get data from the probability μ_{θ_0} , which is the equilibrium probability for the fixed Hölder potential A_{θ_0} .

For convenience, one can assume that the support of Q_0 is the whole space Θ . As an example, Q_0 could be the uniform probability on the simplex Θ (but we do not have to assume that). This could be seen as a form of total ignorance. The initial choice of the probability Q_0 will not be relevant in our setting. Can we identify from the random data, via an iteration procedure, the probability μ_{θ_0} among the others in Θ ?

We will consider once more the log-likelihood point of view.

We assume that the *loss functions* $\ell_n : \Theta \times \Omega \rightarrow \mathbb{R}$, $n \in \mathbb{N}$, are given by

$$\ell_n(\theta, y) = -\log \left(\frac{\mu_\theta(C_n(y))}{\mu_{\theta_0}(C_n(y))} \right) \quad (14)$$

We set

$$\begin{aligned} Z_n(y) &= \int_{\Theta} \exp^{-\ell_n(\theta, y)} dQ_0(\theta) \\ &= \int_{\Theta} \frac{\mu_\theta(C_n(y))}{\mu_{\theta_0}(C_n(y))} dQ_0(\theta) \end{aligned}$$

for each $y \in Y$.

From [20] we get

$$\begin{aligned} \limsup_{n \rightarrow \infty} \frac{1}{n} \log Z_n(y) &= \limsup_{n \rightarrow \infty} \frac{1}{n} \log \int_{\Theta} \frac{\mu_\theta(C_n(y))}{\mu_{\theta_0}(C_n(y))} dQ_0(\theta) \\ &\geq \limsup_{n \rightarrow \infty} \frac{1}{n} \int_{\Theta} \log \frac{\mu_\theta(C_n(y))}{\mu_{\theta_0}(C_n(y))} d\Pi_0(\theta) \\ &= \int_{\Theta} h(\mu_\theta | \mu_{\theta_0}) dQ_0(\theta) > 0 \end{aligned} \quad (15)$$

for μ_{θ_0} -almost every y .

Given $n \in \mathbb{N}$ and $y \in \Omega$, chosen at random according to μ_{θ_0} , we denote by $Q_n = Q_n^y$ the n -posterior probability on Θ , which is given by

$$E \rightarrow Q_n^y(E) = \frac{\int_E \mu_\theta(C_n(y)) dQ_0(\theta)}{\int_{\Theta} \mu_\theta(C_n(y)) dQ_0(\theta)}, \quad (16)$$

on Borel sets $E \subset \Theta$.

In [27] the authors call (16) the n -Gibbs posterior distribution. Given an initial general Q_0 , the proof that

$$\lim_{n \rightarrow \infty} Q_n^y = \delta_{\theta_0},$$

for μ_{θ_0} -a.e. $y \in \Omega$, is a non trivial generalization of the case where just i.i.d. probabilities are considered. Large deviation properties are used in the proof. We will elaborate on the main claim.

Note that for any n, y and E we get

$$Q_n^y(E) = \frac{\int_E \frac{\mu_{\theta}(C_n(y))}{\mu_{\theta_0}(C_n(y))} dQ_0(\theta)}{\int_{\Theta} \frac{\mu_{\theta}(C_n(y))}{\mu_{\theta_0}(C_n(y))} dQ_0(\theta)} = \frac{\int_E \exp^{-\ell_n(\theta, y)} dQ_0(\theta)}{Z_n(y)}. \quad (17)$$

When E is closed and $\theta_0 \notin E$ one can show from (12) that there exists $\alpha > 0$, such that

$$\int_E \frac{\mu_{\theta}(C_n(y))}{\mu_{\theta_0}(C_n(y))} dQ_0(\theta) < \exp^{-\alpha n} Q_0(E). \quad (18)$$

In [20] the following result was shown :

Theorem 2. *Given a Borel set $E \subset \Theta$, the limit*

$$\lim_{n \rightarrow \infty} \frac{1}{n} \log Q_n^y(E) = \lim_{n \rightarrow \infty} \frac{1}{n} \log \frac{\int_E \mu_{\theta}(C_n(y)) dQ_0(\theta)}{\int_{\Theta} \mu_{\theta}(C_n(y)) dQ_0(\theta)} \quad (19)$$

exists for μ_{θ_0} -almost every y . If E is a closed set not containing θ_0 , the above limit is negative. This means $Q_n^y(E)$ goes to zero exponentially fast.

Moreover,

$$\lim_{n \rightarrow \infty} Q_n^y = \delta_{\theta_0}, \quad \text{for } \mu_{\theta_0}\text{-a.e. } y \in \Omega.$$

References

- [1] F. Abramovich and Y. Ritov, *Statistical Theory: A Concise Introduction*, Boca Raton, CRC Press, 2013
- [2] A. Borovkov, *Mathematical Statistics*, Edit. Gordon and Breach
- [3] A. Caticha, *Entropic Physics: Lectures on Probability, Entropy and Statistical Physics*, Lecture Notes, Albany Univ.
<http://dl.icdst.org/pdfs/files1/77964f05542451c01e8e420e975dd664.pdf>

- [4] K-S. Chan and H. Tong, *Chaos: a Statistical Perspective*. Springer-Verlag (2001)
- [5] J-R. Chazottes, R. Floriani and R. Lima, Relative entropy and identification of Gibbs measures in dynamical systems, *J. Statist. Phys.* 90 (1998) no. 3–4, 697–725.
- [6] P. Collet, A. Galves and A. O. Lopes, Maximum Likelihood and Minimum Entropy Estimation of Grammars, *Random and Computational Dynamics*, 3, 241-256, 1995
- [7] M. Denker, Statistical decision procedures and ergodic theory. *Ergodic Theory and Related Topics. Math. Research* **12**, 1982, 35–47.
- [8] M. Denker and G. Keller, On U–statistics and von Mises’ statistics for weakly dependent processes. *Z. Wahrscheinlichkeitsth. verw. Geb.* **64**, 1983, 505–522.
- [9] M. Denker and G. Keller, Rigorous statistical procedures for data from dynamical systems. *J. Stat. Physics* **44**, 1986, 67–93
- [10] M. Denker and M. Gordin, Limit theorems for von Mises statistics of a measure preserving transformation. *Probab. Theory Relat. Fields* 160 1-45 (2014)
- [11] M. Denker, A. O. Lopes and S. R. C. Lopes, Dynamical hypothesis tests and Decision Theory for Gibbs distributions, *Discrete and Continuous Dynamical Systems - Series A - Vol. 43, No. 5*, pp. 1942-1958 (2023)
- [12] M. Denker, M. Kesseböhmer, A. O. Lopes and S. R. C. Lopes, Parametrized Families of Gibbs Measures and their Statistical Inference, *Stoch. and Dynamics*, Vol. 25, No. 02, 2550012 (2025)
- [13] H. H. Ferreira, A. O. Lopes and S. R. C. Lopes, Decision Theory and Large Deviations for Dynamical hypotheses tests: the Neyman-Pearson Lemma, Min-Max and Bayesian Tests, *Journal of Dynamics and Games*, Volume 9, Number 2, 125-150 (2022)
- [14] P. Giulietti, B. Kloeckner, A. O. Lopes and D. Marcon, The calculus of thermodynamical formalism, *Journ. of the European Math Society*, Vol 20, Issue 10, pages 2357–2412 (2018)
- [15] M. J. Karling, S. R. C. Lopes and R. M. de Souza, A Bayesian approach for estimating the parameters of an α -stable distribution, *Journal of Statistical Computation and Simulation*, Vol. 91, issue 9, 1713-1748 (2021)

- [16] M. Karling, S. R. C. Lopes and A. O. Lopes, Explicit Bivariate Rate Functions for Large Deviations in AR(1) and MA(1) Processes with Gaussian Innovations, *Probability, Uncertainty and Quantitative Risk*, Vol. 8, No. 2, 177-212 (2023)
- [17] A. Leucht and M.H. Neumann, Degenerate U - and V -statistics under ergodicity: asymptotics, bootstrap and applications in statistics. *Ann. Inst. Stat. Math.* 65, 349-386 (2013)
- [18] A. O. Lopes, Entropy and Large Deviation, *NonLinearity*, Vol. 3, N. 2, 527-546, 1990.
- [19] A. O. Lopes and J. K. Mengue, The generalized IFS Bayesian method and an associated variational principle covering the classical and dynamical cases, *Dyn. Systems*, Vol. 39, No. 2, 206-230 (2024)
- [20] A. O. Lopes, S. R. C. Lopes and P. Varandas, Bayes posterior convergence for loss functions via almost additive Thermodynamic Formalism, *Journ. of Statis. Physics*, 186:35 (2022)
- [21] A. O. Lopes and J. K. Mengue, On information gain, Kullback-Leibler divergence, entropy production and the involution kernel, *Disc. and Cont. Dyn. Syst. Series A*, Vol. 42, No. 7, 3593-3627 (2022)
- [22] A. O. Lopes and S. R. C. Lopes, Parametric estimation and spectral analysis of piecewise linear maps of the interval, *Adv. in Appl. Probab.* 30 (1998), no. 3, 757-776
- [23] A. O. Lopes and S. R. C. Lopes, Convergence in distribution of the periodogram of chaotic systems, *Stochastic and Dynamics*, vol 2, Issue 4, 609–624 (2002)
- [24] S. R. C. Lopes and M. A. Nunes, Long memory analysis in DNA sequences, *Physica A: Statistical Mechanics and its Applications*, Vol 361, Issue 2, 569-588 (2006)
- [25] B. Olbermann, S. R. C. Lopes and A. O. Lopes, Parameter estimation in Manneville-Pomeau processes, *Probability, Uncertainty and Quantitative Risk*, Vol. 8, No. 2, 213-234 (2023)
- [26] K. McGoff, S. Mukherjee and N.S. Pillai, Statistical inference for dynamical systems: A review. *Statist. Surv.* 9: 209-252 (2015).
- [27] K. McGoff, S. Mukherjee and A. Nobel, Gibbs posterior convergence and Thermodynamic Formalism, *Ann. Appl. Probab.* 32 (2022), no. 1, 461496

- [28] W. Parry and M. Pollicott. Zeta functions and the periodic orbit structure of hyperbolic dynamics, *Astérisque* Vol 187-188 1990
- [29] V. A. Reisen and S. R. C. Lopes, Some simulations and applications of forecasting long-memory time-series models, *Journal of Statistical Planning and Inference*, Vol. 80, Issue 1-2, 269-287 (1999)
- [30] G. Samorodnitsky, *Stochastic Processes and Long Range Dependence*, Springer Verlag (2016)
- [31] A. Zellner, Optimal information processing and Bayes's theorem. *Amer. Statist.* 42 (1988), no. 4, 278-284.
- [32] S. Watanabe, *Mathematical Theory of Bayesian Statistics*, Chapman and Hall/CRC (2018)