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# A SUFFICIENT STATISTIC AND <br> A NOSTANDARD LINEARIZATION IN NONLINEAR REGRESSION 

ANDREJ PÁZMAN

In a nonlinear model $\boldsymbol{y}=\boldsymbol{\eta}(\boldsymbol{\theta})+\boldsymbol{\varepsilon}$ a standard linearization consists in linearizing $\boldsymbol{\eta}(\boldsymbol{\theta})$ at a point $\boldsymbol{\theta}^{*}$, and in computing the M. L. estimate $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{*}\right)$ in the linearized model. We propose to take $\tau(\boldsymbol{y}):=\left(\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right), \ldots, \tau\left(\boldsymbol{y}, \theta^{k}\right)\right)^{\mathrm{T}}$ for some $\boldsymbol{\theta}^{1}, \ldots, \boldsymbol{\theta}^{k}(=$ the sufficient statistic), linearize each $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{i}\right)$ separately, and then to compute the M. L. estimate $\tilde{\boldsymbol{\theta}}(\boldsymbol{y})$. The variable $\tilde{\boldsymbol{\theta}}(\boldsymbol{y})$ has a smaller variance than $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{\boldsymbol{i}}\right)$, and a comparable bias. Further, $\tilde{\boldsymbol{\theta}}(\boldsymbol{y})$ can be used to approximate the posterior density in a Bayesian approach.

The construction of the sufficient statistic has a geometrical background. Possible consequences for nonlinear experimental design are mentioned.

## 1. INTRODUCTION AND THE GEOMETRICAL BACKGROUND

Let us consider the nonlinear regression model with normal errors

$$
\begin{gather*}
\mathbf{y}=\boldsymbol{\eta}(\boldsymbol{\theta})+\boldsymbol{\varepsilon} ; \quad(\boldsymbol{\theta} \in \boldsymbol{\Theta})  \tag{1}\\
\boldsymbol{\varepsilon} \sim \mathscr{N}(\mathbf{0}, \boldsymbol{\Sigma})
\end{gather*}
$$

under standard regularity assumptions: the parameter space $\Theta$ is an open subset of $\mathbb{R}^{m}$, the variance matrix $\boldsymbol{\Sigma}$ is regular, the regression mapping $\boldsymbol{\eta}: \Theta \mapsto \mathbb{R}^{N}(N>m)$ has continuous second order derivatives on $\boldsymbol{\Theta}$, and the vectors $\partial \boldsymbol{\eta}(\boldsymbol{\theta}) / \partial \theta_{1}, \ldots \partial \boldsymbol{\eta}(\boldsymbol{\theta}) / \partial \theta_{\boldsymbol{m}}$ are linearly independent for every $\boldsymbol{\theta} \in \Theta$. The vector $\boldsymbol{y} \in \mathbb{R}^{N}$ is observed, the mapping $\boldsymbol{\eta}$ and the set $\Theta$ are known, $\boldsymbol{\Sigma}$ is either known, or of the form $\boldsymbol{\Sigma}=c \mathbf{W}$ with $c>0$ unknown and $\mathbf{W}$ known. Statistical inference on the unknown vector $\boldsymbol{\theta}$ should be performed.

A well known point estimator in model (1) is the maximum likelihood (= M. L.) estimator

$$
\begin{equation*}
\hat{\boldsymbol{\theta}}:=\hat{\boldsymbol{\theta}}(\boldsymbol{y}):=\underset{\boldsymbol{\theta} \in \boldsymbol{\theta}}{\arg \min }\|\boldsymbol{y}-\boldsymbol{\eta}(\boldsymbol{\theta})\|_{\mathbf{w}}^{2} . \tag{2}
\end{equation*}
$$

Here $\|a\|_{\mathbf{W}}^{2}:=a^{T} \mathbf{W}^{-1} a ; \quad\left(a \in R^{N}\right)$.

In the particular case when model (1) is linear, the statistic $\boldsymbol{y} \in \mathbb{R}^{N} \mapsto \hat{\boldsymbol{\theta}}(\boldsymbol{y})$ is not only a point estimator, it is also a sufficient statistic. If model (1) is nonlinear (more exactly, if the expectation surface of model (1)

$$
\dot{\mathscr{E}}:=\{\boldsymbol{\eta}(\boldsymbol{\theta}): \boldsymbol{\theta} \in \Theta\}
$$

is not a "plane", the statistic $\hat{\boldsymbol{\theta}}(\boldsymbol{y})$ is no more sufficient. Consequently it contains less information about $\boldsymbol{\theta}$ than the sample vector $\boldsymbol{y}$. (For the distributional properties of $\hat{\boldsymbol{\theta}}(\boldsymbol{y})$ cf. e.g. $[4,5])$.

However, it is possible to look for a statistic in model (1) which is a sufficient statistic, and which is somehow related to the M. L. estimator. In particular, we can require that this statistic coincides with $\hat{\boldsymbol{\theta}}(\mathbf{y})$ when model (1) is linear.

In Section 2 we propose such statistics. They have the following geometrical origin:

Consider the expectation surface $\mathscr{E}$. It is an $m$-dimensional surface in the $N$ dimensional sample space $\mathbb{R}^{N}$. According to (2), the point $\boldsymbol{\eta}(\hat{\boldsymbol{\theta}}) \in \mathscr{E}$ is obtained by the $\mathbf{W}$ orthogonal projection of the point $\boldsymbol{y}$ onto $\mathscr{E}$. Consider now for any $\boldsymbol{\theta}^{*} \in \Theta$ the set

$$
T_{\boldsymbol{\theta}^{*}}:=\left\{\boldsymbol{\eta}\left(\boldsymbol{\theta}^{*}\right)+\frac{\partial \boldsymbol{\eta}\left(\boldsymbol{\theta}^{*}\right)}{\partial \boldsymbol{\theta}^{\mathrm{T}}} \mathbf{v}: \mathbf{v} \in \mathbb{R}^{m}\right\} .
$$

Geometrically, $T_{\boldsymbol{\theta}^{*}}$ is the tangent plane to the surface $\mathscr{E}$ at the point $\boldsymbol{\eta}\left(\boldsymbol{\theta}^{*}\right) \in \mathscr{E}$. Statistically, $T_{\boldsymbol{\theta}^{*}}$ is the expectation surface of a linear model which approximates model (1):

$$
\begin{align*}
\boldsymbol{y}-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{*}\right) & =\frac{\partial \boldsymbol{\eta}\left(\theta^{*}\right)}{\partial \boldsymbol{\theta}^{\mathrm{T}}}\left(\boldsymbol{\theta}-\boldsymbol{\theta}^{*}\right)+\boldsymbol{\varepsilon}  \tag{3}\\
\boldsymbol{\varepsilon} & \sim \mathcal{N}(\mathbf{0}, \mathbf{\Sigma})
\end{align*}
$$

The M. L. estimate in this linearized model is

$$
\begin{equation*}
\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{*}\right):=\underset{\boldsymbol{\theta}}{\arg \min }\left\|\boldsymbol{y}-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{*}\right)-\frac{\partial \boldsymbol{\eta}\left(\boldsymbol{\theta}^{*}\right)}{\partial \boldsymbol{\theta}^{\mathrm{T}}}\left(\boldsymbol{\theta}-\boldsymbol{\theta}^{*}\right)\right\|_{\mathbf{w}}^{2} \tag{4}
\end{equation*}
$$

It is the result of the $\mathbf{W}$-orthogonal projection of the point $\boldsymbol{y}$ onto $T_{\boldsymbol{\theta}^{*}}$.
The statistic $\boldsymbol{y} \mapsto \tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{*}\right)$ is sufficient in model (3), however, it is not in model (1). Therefore, we proceed further by considering not one but many (eventually all) tangent planes to $\mathscr{E}$, and by projecting $\mathbf{W}$-orthogonally the sample point $\boldsymbol{y}$ onto all of them. (The reader which is familiar with differential geometry see that we are using the "tangent space" of $\mathscr{E})$. Consequently, instead of one random vector $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{*}\right)$ we consider the set of random vectors

$$
\begin{equation*}
\left\{\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{*}\right) ; \boldsymbol{\theta}^{*} \in D\right\} \tag{5}
\end{equation*}
$$

for some $D \subset \Theta$. Evidently, this is a (vector-valued) random process defined on $D$. This process will be shown to have several pleasant structural properties.
a) It is a Gaussian random process having a covariance function which does not depend on $\boldsymbol{\theta}$.
b) Each component $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{*}\right)$ of this process is related to a linear approximative model.
c) When $D$ is adequately chosen, the mapping

$$
\mathbf{y} \in \mathbb{R}^{N} \mapsto\left\{\tau\left(\mathbf{y}, \boldsymbol{\theta}^{*}\right) ; \boldsymbol{\theta}^{*} \in D\right\}
$$

is a sufficient statistic in model (1).
In Section 3-5 we try to demonstrate that such a process is useful. We restrict our attention to the case of a finite $D=\left\{\boldsymbol{\theta}^{1}, \ldots, \boldsymbol{\theta}^{\boldsymbol{k}}\right\}$, and instead of the process we consider a $k . m$ dimensional random vector $\tau(\boldsymbol{y}):=\left(\tau^{\mathrm{T}}\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right), \ldots, \tau^{\mathrm{T}}\left(\boldsymbol{y}, \boldsymbol{\theta}^{k}\right)\right)^{\mathrm{T}}$. If we linearize each component of $\tau(\boldsymbol{y})$ separately, we obtain a new, nonstandard linearization of model (1) which is more efficient than the standard linearization (3) (see Proposition 2). This allows to obtain an approximative expression for the posterior probability density of $\boldsymbol{\theta}$ (Proposition 3 ). Moreover, using quadratic functions of $\boldsymbol{\tau}(\boldsymbol{y})$ we can discuss some confidence regions for $\boldsymbol{\theta}$, both for the case when $\boldsymbol{\Sigma}$ is known and when $\boldsymbol{\Sigma}=\sigma^{2} \mathbf{I}$ with an unknown $\sigma$.

## 2. SUFFICIENT STATISTICS

As is well known (cf. [1], Chapt. VIII.1.), the M.L. estimate of $\boldsymbol{\theta}$ in the linear model (3) can be expressed in the form

$$
\begin{equation*}
\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{*}\right)=\mathbf{M}^{-1}\left(\boldsymbol{\theta}^{*}\right) \mathbf{F}^{\mathrm{T}}\left(\boldsymbol{\theta}^{*}\right) \mathbf{W}^{-1}\left[\boldsymbol{y}-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{*}\right)\right]+\boldsymbol{\theta}^{*} \tag{6}
\end{equation*}
$$

where

$$
\begin{gathered}
\{\mathbf{F}(\boldsymbol{\theta})\}_{i j}:=\frac{\partial \eta_{i}(\boldsymbol{\theta})}{\partial \theta_{j}} ; \quad(i=1, \ldots, N, \quad j=1, \ldots, m) \\
\mathbf{M}(\boldsymbol{\theta}):=\mathbf{F}^{\mathrm{T}}(\boldsymbol{\theta}) \mathbf{W}^{-1} \mathbf{F}(\boldsymbol{\theta})
\end{gathered}
$$

Consequently, (5) is a Gaussian random process with the mean

$$
\begin{equation*}
\mathbf{m}_{\theta}\left(\boldsymbol{\theta}^{*}\right)=\mathbf{M}^{-1}\left(\boldsymbol{\theta}^{*}\right) \mathbf{F}^{\mathrm{T}}\left(\boldsymbol{\theta}^{*}\right) \mathbf{W}^{-1}\left[\eta(\boldsymbol{\theta})-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{*}\right)\right]+\boldsymbol{\theta}^{*} ; \quad\left(\boldsymbol{\theta}^{*} \in D\right) \tag{7}
\end{equation*}
$$

and with the covariance function $c \mathbf{K}\left(\boldsymbol{\theta}^{*}, \boldsymbol{\theta}^{0}\right)$ where

$$
\mathbf{K}\left(\boldsymbol{\theta}^{*}, \boldsymbol{\theta}^{0}\right)=\mathbf{M}^{-1}\left(\boldsymbol{\theta}^{*}\right) \mathbf{F}^{\mathrm{T}}\left(\boldsymbol{\theta}^{*}\right) \mathbf{W}^{-1} \mathbf{F}\left(\boldsymbol{\theta}^{0}\right) \mathbf{M}^{-1}\left(\boldsymbol{\theta}^{0}\right)
$$

We see that $\mathbf{K}\left(\boldsymbol{\theta}^{*}, \boldsymbol{\theta}^{\boldsymbol{0}}\right)$ does not depend on $\boldsymbol{\theta}(=$ the true value of the parameters).
When the set $D$ is finite, $D=\left\{\boldsymbol{\theta}^{1}, \ldots, \boldsymbol{\theta}^{k}\right\}$, it is better to consider the ( $m . k$ )-dimensional random vector

$$
\tau:=\tau(\boldsymbol{y}):=\left(\begin{array}{c}
\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right)  \tag{8}\\
\vdots \\
\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{k}\right)
\end{array}\right)
$$

instead of the random process (5). Here each component $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{i}\right)$ is defined according to (6).

The mean and the variance matrix of $\tau$ are equal to

$$
\begin{gathered}
\boldsymbol{m}_{\boldsymbol{\theta}}:=\mathrm{E}_{\boldsymbol{\theta}}(\tau)=\left(\boldsymbol{m}_{\boldsymbol{\theta}}^{\mathrm{T}}\left(\boldsymbol{\theta}^{1}\right), \ldots, \boldsymbol{m}_{\boldsymbol{\theta}}^{\mathrm{T}}\left(\boldsymbol{\theta}^{k}\right)\right)^{\mathrm{T}} \\
\operatorname{Var}(\tau)=c \mathbf{S}
\end{gathered}
$$

where

$$
\begin{equation*}
\mathbf{S}:=\binom{\mathbf{K}\left(\boldsymbol{\theta}^{1}, \boldsymbol{\theta}^{1}\right), \ldots, \mathbf{K}\left(\boldsymbol{\theta}^{1}, \boldsymbol{\theta}^{k}\right)}{\mathbf{K}\left(\boldsymbol{\theta}^{k}, \boldsymbol{\theta}^{1}\right), \ldots, \mathbf{K}\left(\boldsymbol{\theta}^{k}, \boldsymbol{\theta}^{k}\right)} \tag{9}
\end{equation*}
$$

If $\mathbf{A}$ is any $r \times s$ matrix, we denote by $\mathscr{M}(\mathbf{A}):=\left\{\mathbf{A} \boldsymbol{u}: \boldsymbol{u} \in \mathbb{R}^{s}\right\}$ the linear subspace of $\mathbb{R}^{r}$ spanned by the columns of $\mathbf{A}$.

Proposition 1. If for every $\theta \in \Theta$

$$
\mathscr{M}[\mathbf{F}(\theta)] \subset \mathscr{M}\left[\left(\mathbf{F}\left(\boldsymbol{\theta}^{1}\right), \ldots, \mathbf{F}\left(\boldsymbol{\theta}^{k}\right)\right)\right]
$$

then the statistic

$$
\mathbf{y} \in \mathbb{R}^{N} \mapsto \tau(\mathbf{y}) \in \mathbb{R}^{m k}
$$

is sufficient in model (1).
Proof. Let $\mathscr{L}$ be the linear manifold in $\mathbb{R}^{N}$ (the "plane") spanned by the set

$$
\bigcup_{\boldsymbol{\theta} \in \boldsymbol{\theta}} T_{\boldsymbol{\theta}} .
$$

Let us define

$$
\mathbf{z}^{\wedge}:=\mathbf{z}^{\wedge}(\mathbf{y}):=\underset{\mathbf{z} \in \mathscr{L}}{\arg \min }\|\mathbf{y}-\mathbf{z}\|_{\mathbf{w}}^{2}
$$

The probability density of $\boldsymbol{y}$ is equal to

$$
\begin{gathered}
f(\mathbf{y} \mid \boldsymbol{\theta})=(2 \pi)^{-N / 2} \operatorname{det}^{-1 / 2}(\boldsymbol{\Sigma}) \exp \left\{-\|\boldsymbol{y}-\eta(\boldsymbol{\theta})\|_{\mathbf{W}}^{2} /(2 c)\right\} \approx \\
\approx \exp \left\{-\left\|\boldsymbol{y}-\mathbf{z}^{\wedge}(\boldsymbol{y})\right\|_{\mathbf{W}}^{2} /(2 c)\right\} \exp \left\{-\left\|\mathbf{z}^{\wedge}-\boldsymbol{\eta}(\boldsymbol{\theta})\right\|_{\mathbf{W}}^{2} /(2 c\} ; \quad(\boldsymbol{\theta} \in \boldsymbol{\Theta})\right.
\end{gathered}
$$

Hence, according to the factorisation theorem (cf. [1], Chapt. XV. 5.), the statistic $\boldsymbol{z}^{\wedge}(\boldsymbol{y})$ is sufficient in model (1).

Denote by

$$
\mathbf{P}_{i}:=\mathbf{F}\left(\boldsymbol{\theta}^{i}\right) \mathbf{M}^{-1}\left(\boldsymbol{\theta}^{i}\right) \mathbf{F}^{\mathrm{T}}\left(\boldsymbol{\theta}^{i}\right) \mathbf{W}^{-1}
$$

the $\mathbf{W}$-orthogonal projector onto $\mathscr{M}\left[\mathbf{F}\left(\boldsymbol{\theta}^{i}\right)\right]$. The mapping $\mathbf{z} \mapsto\left(\mathbf{P}_{1}\left(\mathbf{z}-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{1}\right)\right), \ldots\right.$ $\ldots, \mathbf{P}_{k}\left(\boldsymbol{z}-\eta\left(\boldsymbol{\theta}^{k}\right)\right)$ is one-to one on $\mathscr{L}$. Indeed, take $\boldsymbol{z}, \mathbf{z}^{*} \in \mathscr{L}$ such that

$$
\mathbf{P}_{i}\left(\mathbf{z}-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{i}\right)\right)=\mathbf{P}_{i}\left(\mathbf{z}^{*}-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{i}\right)\right) ; \quad(i=1, \ldots, k)
$$

Multiplying by $\mathbf{F}^{\mathrm{T}}\left(\boldsymbol{\theta}^{i}\right) \mathbf{W}^{-1}$ from the left, we obtain

$$
\mathbf{F}^{\mathrm{T}}\left(\boldsymbol{\theta}^{i}\right) \mathbf{W}^{-1}\left(\mathbf{z}-\mathbf{z}^{*}\right)=0 ; \quad(i=1, \ldots, k)
$$

i.e. $\left(\boldsymbol{z}-\mathbf{z}^{*}\right)$ is $\mathbf{W}$-orthogonal to $\mathscr{M}\left[\left(\mathbf{F}\left(\boldsymbol{\theta}^{1}\right), \ldots, \mathbf{F}\left(\boldsymbol{\theta}^{k}\right)\right)\right]$, hence to $\mathscr{L}$. Consequently, $\left(\mathbf{z}-\mathbf{z}^{*}\right)^{\mathrm{T}} \mathbf{W}^{-1}\left(\mathbf{z}-\mathbf{z}^{*}\right)=0$ hence, $\mathbf{z}-\mathbf{z}^{*}$.

It follows that $\boldsymbol{y} \in \mathbb{R}^{N} \mapsto\left(\mathbf{P}_{\mathbf{1}}\left(\mathbf{z}^{\wedge}(\boldsymbol{y})-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{\mathbf{1}}\right)\right), \ldots, \mathbf{P}_{k}\left(\mathbf{z}^{\wedge}(\boldsymbol{y})-\eta\left(\boldsymbol{\theta}^{k}\right)\right)\right)$ is a sufficient statistic in model (1).

Since $\boldsymbol{z}^{\wedge}(\boldsymbol{y})$ is the $\mathbf{W}$-orthogonal projection of $\boldsymbol{y}$ onto $\mathscr{L}$ we have

$$
\mathbf{P}_{i}\left(\mathbf{z}^{\wedge}(\boldsymbol{y})-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{i}\right)\right)=\mathbf{P}_{i}\left(\boldsymbol{y}-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{i}\right)\right) ; \quad(i=1, \ldots, k)
$$

Further, the equality

$$
\mathbf{F}\left(\boldsymbol{\theta}^{i}\right) \tau\left(\mathbf{y}, \boldsymbol{\theta}^{i}\right)=\mathbf{P}_{i}\left(\boldsymbol{y}-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{i}\right)\right)+\mathbf{F}\left(\boldsymbol{\theta}^{i}\right) \boldsymbol{\theta}^{i}
$$

which follows from Eq. (6), specifies $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{i}\right)$ uniquely, since $\mathbf{F}\left(\boldsymbol{\theta}^{\boldsymbol{i}}\right)$ is of full rank. Consequently the mapping $\tau(\boldsymbol{y}) \mapsto\left(\mathbf{P}_{\mathbf{1}}\left(\mathbf{z}^{\wedge}(\boldsymbol{y})-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{1}\right)\right), \ldots, \mathbf{P}_{\boldsymbol{k}}\left(\mathbf{z}^{\wedge}(\boldsymbol{y})-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{\boldsymbol{k}}\right)\right)\right.$ is one-to-one. It follows that $\tau(\boldsymbol{y})$ is a sufficient statistic in model (1).

Corollary 1. If $D \subset \Theta$ is such that

$$
\mathscr{M}[\mathbf{F}(\theta)] \subset \mathscr{M}\left[\left(\mathbf{F}\left(\theta^{1}\right), \ldots, \mathbf{F}\left(\theta^{k}\right)\right)\right] ; \quad(\theta \in \Theta)
$$

for some finite set $\left\{\boldsymbol{\theta}^{1}, \ldots, \boldsymbol{\theta}^{\boldsymbol{k}}\right\} \subset D$, then

$$
\boldsymbol{y} \in \mathbb{R}^{N} \mapsto\left\{\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{*}\right) ; \quad \boldsymbol{\theta}^{*} \in D\right\}
$$

is sufficient. Particularly

$$
\boldsymbol{y} \in \mathbb{R}^{N} \mapsto\left\{\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{*}\right) ; \boldsymbol{\theta}^{*} \in \boldsymbol{\Theta}\right\}
$$

is always sufficient.
Corollary 2. Let $\pi(\boldsymbol{\theta})$ be a probability density on $\Theta$ (the prior density) such that

$$
\mathscr{M}[\mathbf{F}(\boldsymbol{\theta})] \subset \mathscr{M}\left[\left(\mathbf{F}\left(\boldsymbol{\theta}^{1}\right), \ldots, \mathbf{F}\left(\boldsymbol{\theta}^{k}\right)\right)\right] ; \quad(\boldsymbol{\theta} \in \operatorname{supp}(\pi)) .
$$

Then

$$
\pi(\boldsymbol{\theta} \mid \tau(\mathbf{y}))=\pi(\boldsymbol{\theta} \mid \boldsymbol{y})
$$

where $\pi(\boldsymbol{\theta} \mid \boldsymbol{u})$ denotes the posterior density of $\boldsymbol{\theta}$ given $\boldsymbol{u}$.
Proof. As in Proposition 1, we can prove that $\tau(\boldsymbol{y})$ is sufficient in the model

$$
\boldsymbol{y}=\boldsymbol{\eta}(\boldsymbol{\theta})+\boldsymbol{\varepsilon} ; \quad(\boldsymbol{\theta} \in \operatorname{supp}(\pi)) .
$$

Hence $f(\boldsymbol{y} \mid \theta)$ can be factorized, i.e. we can write

$$
f(\mathbf{y} \mid \boldsymbol{\theta})=h(\mathbf{y}) g(\tau(\mathbf{y}), \boldsymbol{\theta})
$$

for some functions $h$ and $g$. It follows that

$$
\pi(\boldsymbol{\theta} \mid \mathbf{y})=\frac{f(\mathbf{y} \mid \boldsymbol{\theta}) \pi(\boldsymbol{\theta})}{\int_{\operatorname{supp}(\pi)} f(\mathbf{y} \mid \boldsymbol{t}) \pi(\mathbf{t}) \mathrm{d} \mathbf{t}}=\frac{g(\tau(\mathbf{y}), \boldsymbol{\theta}) \pi(\boldsymbol{\theta})}{\int_{\operatorname{supp}(\pi)} g(\tau(\mathbf{y}), \boldsymbol{t}) \pi(\mathbf{t}) \mathrm{d} \boldsymbol{t}}
$$

Hence $\boldsymbol{y} \in \mathbb{R}^{N} \mapsto \pi(\boldsymbol{\theta} \mid \boldsymbol{y})$ is a function of $\tau(\boldsymbol{y})$. According to the definition of conditional distributions (cf. [7], Chapt. V. 1.) it means that $\pi(\boldsymbol{\theta} \mid \boldsymbol{y})=\pi(\boldsymbol{\theta} \mid \tau(\boldsymbol{y}))$.

## 3. A NONSTANDARD LINEARIZATION

Let us consider the random vector $\boldsymbol{\tau}(\mathbf{y})$ (the sufficient statistic) defined in Eq. (8). We have

$$
\begin{equation*}
\tau(\boldsymbol{y}) \sim \mathscr{N}\left(\boldsymbol{m}_{\boldsymbol{\theta}}, c \mathbf{S}\right) ; \quad(\boldsymbol{\theta} \in \boldsymbol{\Theta}) \tag{10}
\end{equation*}
$$

where $\boldsymbol{m}_{\boldsymbol{\theta}}$ and $\mathbf{S}$ are given by (9).
Instead of taking the linearization (3) we propose to linearize (10), i.e. to take
approximatively

$$
\begin{equation*}
\tau(\boldsymbol{y}) \sim \mathscr{N}(\mathbf{J} \boldsymbol{\theta}, c \mathbf{S}) ; \quad\left(\boldsymbol{\theta} \in \mathbb{R}^{m}\right) \tag{11}
\end{equation*}
$$

where

$$
\mathbf{J}:=(\mathbf{I}, \ldots, \mathbf{I})^{\mathrm{T}}
$$

and $\mathbf{I}$ is the $m \times m$ identity matrix. The linearization (11) is the linearization (3) applied separately to each component $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{i}\right)$ of the vector $\tau(\boldsymbol{y})$.

To compare the standard linearization (3) with (11) take for $\boldsymbol{\theta}^{*}$ any point of the set $\left\{\boldsymbol{\theta}^{1}, \ldots, \boldsymbol{\theta}^{k}\right\}$, say $\boldsymbol{\theta}^{*}=\boldsymbol{\theta}^{1}$. Then consider the BLUE-s ( $=$ best linear unbiased estimates) of $\boldsymbol{\theta}$ in both models. The BLUE in model (3) is equal to $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right)$, and is expressed in Eq. (6). Althoug the matrix $\mathbf{S}$ is singular (in general), and $\mathscr{M}(\mathbf{J}) \nsubseteq \mathscr{M}(\mathbf{S})$, the vector $\boldsymbol{\theta}$ can be estimated without bias in model (11), say by the estimate

$$
\frac{1}{k} \mathbf{J}^{\mathrm{T}} \tau(\mathbf{y})
$$

Hence the BLUE exists also in model (11). Let us denote it by $\tilde{\boldsymbol{\theta}}(\mathbf{y})$. We refere to [3], Theorems 5.2.2 and 5.2.5 for explicit expressions for $\tilde{\boldsymbol{\theta}}(\mathbf{y})$ and $\operatorname{Var} \tilde{\boldsymbol{\theta}}(\mathbf{y})$. We have

$$
\tilde{\boldsymbol{\theta}}(\mathbf{y})=\mathbf{Q} \tau(\mathbf{y}), \quad \operatorname{Var} \tilde{\boldsymbol{\theta}}(\mathbf{y})=c \mathbf{V}
$$

where

$$
\begin{align*}
& \mathbf{Q}:=\left[\mathbf{J}^{\mathrm{T}}\left(\mathbf{S}+\mathbf{J} \mathbf{J}^{\mathrm{T}}\right)^{-} \mathbf{J}\right]^{-1} \mathbf{J}^{\mathrm{T}}\left(\mathbf{S}+\mathbf{J} \mathbf{J}^{\mathrm{T}}\right)^{-}  \tag{12}\\
& \mathbf{V}:=\left[\mathbf{J}^{\mathrm{T}}\left(\mathbf{S}+\mathbf{J} \mathbf{J}^{\mathrm{T}}\right)^{-} \mathbf{J}\right]^{-1}-\mathbf{I}
\end{align*}
$$

We note that $\mathbf{J}^{\mathrm{T}}\left(\mathbf{S}+\mathbf{J} \mathbf{J}^{\mathbf{T}}\right)^{-} \mathbf{J}$ is nonsingular, since $\mathbf{J}$ is of full rank and $\mathscr{M}[\mathbf{J}]=$ $=\mathscr{M}\left[\mathbf{J} \mathbf{J}^{\mathrm{T}}\right] \subset \mathscr{M}\left[\mathbf{S}+\mathbf{J J}^{\mathrm{T}}\right]$. In the particular case that $\mathbf{S}$ is regular, we have simpler formulae

$$
\begin{align*}
& \mathbf{Q}=\left(\mathbf{J}^{\mathrm{T}} \mathbf{S}^{-1} \mathbf{J}\right)^{-1} \mathbf{J}^{\mathrm{T}} \mathbf{S}^{-1}  \tag{13}\\
& \mathbf{V}=\left(\mathbf{J}^{\mathrm{T}} \mathbf{S}^{-1} \mathbf{J}\right)^{-1}
\end{align*}
$$

Hence in the linearized model (11) we have

$$
\begin{equation*}
\tilde{\boldsymbol{\theta}}(\mathbf{y}) \sim \mathscr{N}(\boldsymbol{\theta}, c \mathbf{V}) ; \quad\left(\boldsymbol{\theta} \in \mathbb{R}^{m}\right) \tag{14}
\end{equation*}
$$

but in the linearized model (3) we have

$$
\begin{equation*}
\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right) \sim \mathscr{N}\left(\boldsymbol{\theta}, c \mathbf{M}^{-1}\left(\boldsymbol{\theta}^{1}\right)\right) ; \quad\left(\boldsymbol{\theta} \in \mathbb{R}^{m}\right) \tag{15}
\end{equation*}
$$

To compare what linearization is better, we shall compare the exact distributions of $\tilde{\boldsymbol{\theta}}(\boldsymbol{y})$ and $\boldsymbol{\tau}\left(\boldsymbol{y}, \boldsymbol{\theta}^{\boldsymbol{v}}\right)$.

Proposition 2. The random vectors $\tilde{\boldsymbol{\theta}}(\mathbf{y})$ and $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right)$ are exactly distributed according to

$$
\begin{gather*}
\tilde{\boldsymbol{\theta}}(\mathbf{y}) \sim \mathscr{N}\left(\mathbf{Q} \boldsymbol{m}_{\boldsymbol{\theta}}, c \mathbf{V}\right) ; \quad(\boldsymbol{\theta} \in \boldsymbol{\Theta})  \tag{16}\\
\tau\left(\mathbf{y}, \boldsymbol{\theta}^{1}\right) \sim \mathscr{N}\left(\boldsymbol{m}_{\boldsymbol{\theta}}\left(\boldsymbol{\theta}^{1}\right), c \mathbf{M}^{-1}\left(\boldsymbol{\theta}^{1}\right)\right) ; \quad(\boldsymbol{\theta} \in \boldsymbol{\Theta}) \tag{17}
\end{gather*}
$$

The vectors expressing the bias

$$
\mathbf{Q} \boldsymbol{m}_{\theta}-\theta
$$

and the bias

$$
m_{\theta}\left(\theta^{1}\right)-\theta
$$

are of the sample order of magnitude. The estimator $\tilde{\boldsymbol{\theta}}(\boldsymbol{y})$ is more efficient since the matrix $\operatorname{Var}\left[\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right)\right]-\operatorname{Var}[\tilde{\boldsymbol{\theta}}(\boldsymbol{y})]$ is positive semidefinite.

Proof. Both variables $\tilde{\boldsymbol{\theta}}(\boldsymbol{y})$ and $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{\mathbf{1}}\right)$ are linear in $\boldsymbol{y}$, hence they are normally distributed. The mean and the variance of $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right)$ is given in Eq. (7). The mean and the variance of $\tilde{\boldsymbol{\theta}}(\boldsymbol{y})$ follow from Eq. (12) and from the mean and the variance of $\tau(y)$ in Eq. (9).

The bias of $\tilde{\theta}(\mathbf{y})$ is

$$
\mathbf{Q} m_{\theta}-\theta=\mathbf{Q}\left(\begin{array}{c}
m_{\theta}\left(\theta^{1}\right) \\
\vdots \\
m_{\theta}\left(\theta^{k}\right)
\end{array}\right)-\theta
$$

The bias of $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right)$ can be written in the form

$$
m_{\theta}\left(\theta^{1}\right)-\theta=Q \mathbf{J} m_{\theta}\left(\theta^{1}\right)-\theta=\mathbf{Q}\left(\begin{array}{c}
m_{\theta}\left(\theta^{1}\right) \\
\vdots \\
m_{\theta}\left(\theta^{1}\right)
\end{array}\right)-\theta
$$

since $\mathbf{Q} \mathbf{J}=\mathbf{I}$, according to (12) and (13). Thus if $\boldsymbol{m}_{\boldsymbol{\theta}}\left(\boldsymbol{\theta}^{\boldsymbol{i}}\right)-\boldsymbol{\theta}$ is of the same order for every $i=1, \ldots, k$, then $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right)$ and $\tilde{\boldsymbol{\theta}}(\boldsymbol{y})$ have the bias of the same order as well.

The random variable $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right)$ can be written in the form

$$
\tau\left(\mathbf{y}, \boldsymbol{\theta}^{1}\right)=(\mathbf{I}, \mathbf{0}, \ldots, \mathbf{0}) \tau(\mathbf{y})
$$

hence it is a linear unbiassed estimator of $\boldsymbol{\theta}$ in model (11). Since $\tilde{\boldsymbol{\theta}}(\boldsymbol{y})$ is the BLUE in the same model, it follows that $\operatorname{Var}\left[\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right)\right]-\operatorname{Var}[\tilde{\boldsymbol{\theta}}(\boldsymbol{y})]$ is positive semidefinite.

Note 1. According to Eq. (7) we can write the bias in the form

$$
m_{\theta}\left(\theta^{1}\right)-\theta:=r\left(\theta, \theta^{1}\right), \quad J m_{\theta}-\theta=\mathbf{J}\left(\begin{array}{c}
r\left(\theta, \theta^{1}\right) \\
\vdots \\
r\left(\theta, \theta^{k}\right)
\end{array}\right)
$$

where from the Taylor formula for $\boldsymbol{\eta}(\boldsymbol{\theta})$ at $\boldsymbol{\theta}^{\boldsymbol{i}}$ we obtain

$$
\begin{gather*}
\boldsymbol{r}\left(\boldsymbol{\theta}, \boldsymbol{\theta}^{i}\right):=\mathbf{M}^{-1}\left(\boldsymbol{\theta}^{i}\right) \mathbf{F}^{\mathrm{T}}\left(\boldsymbol{\theta}^{i}\right) \mathbf{W}^{-1}\left[\eta(\boldsymbol{\theta})-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{i}\right)\right]+\boldsymbol{\theta}^{i}-\boldsymbol{\theta}=  \tag{18}\\
=\frac{1}{2} \mathbf{M}^{-1}\left(\boldsymbol{\theta}^{i}\right) \mathbf{F}^{\mathrm{T}}\left(\boldsymbol{\theta}^{i}\right) \mathbf{W}^{-1}\left[( \boldsymbol { \theta } - \boldsymbol { \theta } ^ { i } ) ^ { \mathrm { T } } \left[\frac{\partial^{2} \boldsymbol{\eta}(\boldsymbol{\theta})}{\left.\left.\left.\partial \boldsymbol{\theta} \frac{\boldsymbol{\theta}^{\mathrm{T}}}{}\right|_{\lambda \boldsymbol{\theta}+(1-\lambda) \boldsymbol{\theta}^{i}}\right]\left(\boldsymbol{\theta}-\boldsymbol{\theta}^{i}\right)\right]}\right.\right.
\end{gather*}
$$

for some number $\lambda \in(0,1)$ depending on $\boldsymbol{\theta}$ and on $\boldsymbol{\theta}^{\boldsymbol{i}}$.
The expression for $\boldsymbol{r}\left(\boldsymbol{\theta}, \boldsymbol{\theta}^{\boldsymbol{i}}\right)$ is small either if $\left[\boldsymbol{\theta}-\boldsymbol{\theta}^{\boldsymbol{i}}\right]^{\mathrm{T}} \mathbf{M}\left(\boldsymbol{\theta}^{\boldsymbol{i}}\right)\left[\boldsymbol{\theta}-\boldsymbol{\theta}^{\boldsymbol{i}}\right]$ is small or if model (1) is not too much curved, since

$$
\sup \left\{\left.\left\|\boldsymbol{v}^{\mathrm{T}} \frac{\partial^{2} \boldsymbol{\eta}(\boldsymbol{\theta})}{\partial \boldsymbol{\theta} \partial \boldsymbol{\theta}^{\mathrm{T}}}\right\|_{\mathbf{w}} \right\rvert\, \mathbf{v}^{\mathrm{T}} \mathbf{M}\left(\boldsymbol{\theta}^{1}\right) \mathbf{v} ; \quad 0 \neq \mathbf{v} \in \mathbb{R}^{m}\right\}
$$

is related to the curvatures of Bates and Watts [2] in model (1). We used here the notation

$$
\mathbf{v}^{\mathrm{T}} \frac{\partial^{2} \boldsymbol{\eta}(\boldsymbol{\theta})}{\partial \boldsymbol{\theta} \partial \boldsymbol{\theta}^{\mathrm{T}}} \boldsymbol{v}:=\sum_{i j} v_{i} \frac{\partial^{2} \boldsymbol{\eta}(\boldsymbol{\theta})}{\partial \theta_{i} \partial \theta_{j}} v_{j} .
$$

It is important to note here that $\mathrm{E}_{\theta}[\tilde{\theta}(\boldsymbol{y})]=\mathbf{Q} \boldsymbol{m}_{\theta}$ is a "mixture" of the means of $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{1}\right), \ldots, \tau\left(\boldsymbol{y}, \tilde{\boldsymbol{\theta}}^{k}\right)$. In some cases the "mixture" is such that the bias of $\tilde{\boldsymbol{\theta}}(\boldsymbol{y})$ is much smaller than the bias of every $\tau\left(\boldsymbol{y}, \boldsymbol{\theta}^{i}\right)$. This depends on the choice off $\boldsymbol{\theta}^{1}, \ldots, \boldsymbol{\theta}^{k}$.

Note 2. When $\boldsymbol{\tau}(\boldsymbol{y})$ is a sufficient statistic (Proposition 1) we arrive to $\tilde{\boldsymbol{\theta}}(\boldsymbol{y})$ according to the scheme

$$
\begin{gathered}
(1) \mapsto \text { sufficient statistic } \tau \mapsto(10) \mapsto \text { linearization of } \tau \mapsto(11) \quad \mapsto \\
\mapsto \text { sufficient statistic } \tilde{\boldsymbol{\theta}} \mapsto(14)
\end{gathered}
$$

Example 1. We shall consider the simple nonlinear model $\boldsymbol{y}=\boldsymbol{\eta}(\boldsymbol{\theta})+\boldsymbol{\varepsilon}$ with $N=2, m=1, \boldsymbol{\Theta}=(0, \pi), \boldsymbol{\eta}(\boldsymbol{\theta})=(\cos \theta, \sin \theta)^{\mathrm{T}}, \boldsymbol{\Sigma}=\mathbf{W}=\mathbf{I}$. (The expectation surface is a halfcircle). In this case we have $\mathbf{F}^{\mathrm{T}}(\theta)=(-\sin \theta, \cos \theta), \mathbf{M}(\theta)=1$; $(\theta \in \Theta)$. To construct $\tau(\boldsymbol{y})$ take two points $\theta^{1}=\theta^{*}-\delta, \theta^{2}=\theta^{*}+\delta$ for some fixed $\delta>0, \theta^{*} \in \Theta$. By simple computations we obtain

$$
\begin{aligned}
& \tau(\mathbf{y})=\binom{-y_{1} \sin \left(\theta^{*}-\delta\right)+y_{2} \cos \left(\theta^{*}-\delta\right)+\theta^{*}-\delta}{-y_{1} \sin \left(\theta^{*}+\delta\right)+y_{2} \cos \left(\theta^{*}+\delta\right)+\theta^{*}+\delta} \\
& \mathbf{S}=\left(\begin{array}{cc}
1, & \cos (2 \delta) \\
\cos (2 \delta), & 1
\end{array}\right), \quad \mathbf{S}^{-1}=\left(\begin{array}{cc}
1, & -\cos (2 \delta) \\
-\cos (2 \delta), & 1
\end{array}\right) / \sin ^{2}(2 \delta) \\
& \mathbf{V}=\left(\mathbf{J}^{\mathrm{T}} \mathbf{S}^{-1} \mathbf{J}\right)^{-1}=(1+\cos 2 \delta) / 2=\cos ^{2} \delta<1 \\
& \tilde{\theta}(\mathbf{y})=\left(\mathbf{J}^{\mathrm{T}} \mathbf{S}^{-1} \mathbf{J}\right)^{-1} \mathbf{J}^{\mathrm{T}} \mathbf{S}^{-1} \tau(\mathbf{y})=\cos \delta\left[-y_{1} \sin \theta^{*}+y_{2} \cos \theta^{*}\right]+\theta^{*}
\end{aligned}
$$

When $\delta \mapsto 0$ we obtain $\tau\left(\mathbf{y}, \theta^{*}\right)$ :

$$
\tau\left(\mathbf{y}, \theta^{*}\right)=\left[-y_{1} \sin \theta^{*}+y_{2} \cos \theta^{*}\right]+\theta^{*}
$$

Further

$$
\begin{aligned}
& \mathrm{E}_{\theta}[\tilde{\theta}(\mathbf{y})]-\theta=\cos \delta \sin \left(\theta-\theta^{*}\right)+\left(\theta^{*}-\theta\right) \\
& \mathrm{E}_{\theta}\left[\tau\left(\mathbf{y}, \theta^{*}\right)\right]-\theta=\sin \left(\theta-\theta^{*}\right)+\left(\theta^{*}-\theta\right)
\end{aligned}
$$

Hence for $\delta$ not very large, the bias of $\tilde{\theta}(\boldsymbol{y})$ and of $\tau\left(\boldsymbol{y}, \theta^{*}\right)$ is approximatively the same. The mean square error of $\tilde{\theta}(\mathbf{y})$ is equal to

$$
\mathrm{E}_{\theta}[\tilde{\theta}(\boldsymbol{y})-\theta]^{2}=\cos ^{2} \delta+\left[\cos \delta \sin \left(\theta-\theta^{*}\right)+\left(\theta^{*}-\theta\right)\right]^{2}:=\psi(\delta)
$$

We have

$$
\frac{\mathrm{d} \psi}{\mathrm{~d} \delta}=-2 \cos \delta \sin \delta\left[1+\sin ^{2}\left(\theta-\theta^{*}\right)\right]-2 \sin \delta\left[\sin \left(\theta-\theta^{*}\right)\right]\left(\theta^{*}-\theta\right)
$$

Hence $\mathrm{d} \psi /\left.\mathrm{d} \delta\right|_{\delta=0}=0$. Further

$$
\begin{aligned}
\left.\frac{\mathrm{d}^{2} \psi}{\mathrm{~d} \delta^{2}}\right|_{\delta=0}= & -\left.2\left[-\sin ^{2} \delta+\cos ^{2} \delta\right]\left[1+\sin ^{2}\left(\theta-\theta^{*}\right)\right]\right|_{\delta=0}- \\
& -\left.2 \cos \delta\left[\sin \left(\theta-\theta^{*}\right)\right]\left(\theta^{*}-\theta\right)\right|_{\delta=0}= \\
= & -2\left[1+\sin ^{2}\left(\theta-\theta^{*}\right)\right]+2\left(\theta-\theta^{*}\right) \sin \left(\theta-\theta^{*}\right)
\end{aligned}
$$

If $\theta>\theta^{*}$ then

$$
\left.\frac{\mathrm{d}^{2} \psi}{\mathrm{~d} \delta^{2}}\right|_{\delta=0} \leqq-2+2\left[\theta-\theta^{*}-\sin \left(\theta-\theta^{*}\right)\right]
$$

If $\theta<\theta^{*}$ then

$$
\left.\left.\frac{\mathrm{d}^{2} \psi}{\mathrm{~d} \delta^{2}}\right|_{\delta=0} \leqq-2+2\left[\theta^{*}-\theta-\sin \left(\theta^{*}-\theta\right)\right)\right]
$$

Hence, if $\theta^{*}$ is so near to $\theta$ that

$$
\left|\theta-\theta^{*}\right|-\sin \left(\left|\theta-\theta^{*}\right|\right)<1
$$

then $\mathrm{E}_{\theta}[\tilde{\theta}(\mathbf{y})-\theta]^{2}$ attains its maximum at $\delta=0$. Consequently

$$
\mathrm{E}_{\theta}[\tilde{\theta}(\mathbf{y})-\theta]^{2}<\mathrm{E}_{\theta}\left[\tau\left(\mathbf{y}, \theta^{*}\right)-\theta\right]^{2}
$$

## 4. THE POSTERIOR PROBABILITY DENSITY OF $\theta$

Consider a normal prior density $\pi(\boldsymbol{\theta})$ in model (1),

$$
\pi(\boldsymbol{\theta})=(2 \pi)^{-m / 2} \operatorname{det}^{-1 / 2}(\mathbf{H}) \exp \left\{-\frac{1}{2}\left(\boldsymbol{\theta}-\boldsymbol{\theta}^{0}\right)^{\mathrm{T}} \mathbf{H}^{-1}\left(\boldsymbol{\theta}-\boldsymbol{\theta}^{0}\right)\right\}
$$

where $\mathbf{H}$ is a given matrix and $\boldsymbol{\theta}^{0} \in \Theta$ is a given vector. Denote by $\pi(\boldsymbol{\theta} \mid \boldsymbol{y})$ the corresponding posterior density. If $\tau(\boldsymbol{y})$ is a sufficient static (Corollary 2 to Proposition 1) then

$$
\pi(\boldsymbol{\theta} \mid \mathbf{y})=\pi(\boldsymbol{\theta} \mid \tau(\mathbf{y}))
$$

This is not a normal density. However, using the linearization described in Section 3, we can write approximatively

$$
\pi(\boldsymbol{\theta} \mid \boldsymbol{y}) \doteq \pi_{\operatorname{lin}}(\boldsymbol{\theta} \mid \tilde{\boldsymbol{\theta}}(\boldsymbol{y}))
$$

where $\tilde{\theta}(\mathbf{y})$ is supposed to be distributed according to Eq. (14).
Proposition 3. $\pi_{\text {lin }}(\boldsymbol{\theta} \mid \tilde{\boldsymbol{\theta}}(\mathbf{y}))$ is a normal probability density with the mean equal to

$$
\begin{equation*}
\boldsymbol{\theta}^{0}+\mathbf{H}(c \mathbf{V}+\mathbf{H})^{-1}\left(\tilde{\boldsymbol{\theta}}(\mathbf{y})-\boldsymbol{\theta}^{0}\right) \tag{19}
\end{equation*}
$$

and with the variance matrix equal to

$$
\begin{equation*}
\mathbf{H}-\mathbf{H}(c \mathbf{V}+\mathbf{H})^{-1} \mathbf{H} \tag{20}
\end{equation*}
$$

where $\tilde{\boldsymbol{\theta}}(\mathbf{y})$ and $\mathbf{V}=\operatorname{Var} \tilde{\boldsymbol{\theta}}(\boldsymbol{y})$ are defined by Eqs. (12) resp. (13).
Proof. Denote by $h_{\text {lin }}(\tilde{\boldsymbol{\theta}} \mid \boldsymbol{\theta})$ the probability density of $\tilde{\boldsymbol{\theta}}$ corresponding to Eq. (14).

Consider the vector
as a random vector with the joint density $\pi(\boldsymbol{\theta}) h_{\text {lin }}(\tilde{\theta} \mid \boldsymbol{\theta})$ and denote by $E(\cdot)$ the operator of taking the mean with respect to this density. By simple computations we obtain

$$
\begin{aligned}
& \mathrm{E}(\boldsymbol{\theta})=\boldsymbol{\theta}^{0} \\
& \mathrm{E}(\tilde{\boldsymbol{\theta}})=\mathrm{E}\left[\mathrm{E}_{\boldsymbol{\theta}}(\tilde{\boldsymbol{\theta}})\right]=\boldsymbol{\theta}^{0} \\
& \mathrm{E}\left[\left(\boldsymbol{\theta}-\boldsymbol{\theta}^{0}\right)\left(\theta-\boldsymbol{\theta}^{0}\right)^{\mathrm{T}}\right]=\mathbf{H} \\
& \mathrm{E}\left[\left(\boldsymbol{\theta}-\boldsymbol{\theta}^{0}\right)\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}^{0}\right)^{\mathrm{T}}=\mathrm{E}\left[\left(\boldsymbol{\theta}-\boldsymbol{\theta}^{0}\right) \mathrm{E}_{\boldsymbol{\theta}}\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}^{0}\right)^{\mathrm{T}}\right]=\mathbf{H}\right. \\
& \mathrm{E}\left[\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}^{0}\right)\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}^{0}\right)^{\mathrm{T}}\right]=\mathrm{E}\left[\mathrm{E}_{\boldsymbol{\theta}}\left[\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}^{0}\right)\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}^{0}\right)^{\mathrm{T}}\right]\right]=c \mathbf{V}+\mathbf{H} .
\end{aligned}
$$

Hence

$$
\binom{\tilde{\boldsymbol{\theta}}}{\boldsymbol{\theta}} \sim \mathscr{N}\left(\binom{\boldsymbol{\theta}^{0}}{\boldsymbol{\theta}^{0}},\binom{c \mathbf{V}+\mathbf{H}, \mathbf{H}}{\mathbf{H}, \mathbf{H}}\right) .
$$

According to [6], Chapt. 8. a $2,(\mathrm{~V})$, the conditional density of $\boldsymbol{\theta}$ given $\tilde{\boldsymbol{\theta}}$ is normal with the mean

$$
\boldsymbol{\theta}^{0}+\mathbf{H}(c \mathbf{V}+\mathbf{H})^{-1}\left(\tilde{\boldsymbol{\theta}}-\boldsymbol{\theta}^{0}\right)
$$

and with the variance

$$
\mathbf{H}-\mathbf{H}(c \mathbf{V}+\mathbf{H})^{-1} \mathbf{H} .
$$

Note. The statistic $\tilde{\theta}(\boldsymbol{y})$ is sufficient in the linearized model (11). Therefore we can write (compare with Corollary 2)

$$
\pi_{\mathrm{lin}}(\boldsymbol{\theta} \mid \tilde{\theta}(\mathbf{y}))=\pi_{\operatorname{lin}}(\boldsymbol{\theta} \mid \tau(\mathbf{y}))
$$

On the other hand, the exact posterior density is

$$
\pi(\boldsymbol{\theta} \mid \boldsymbol{y})=\pi(\boldsymbol{\theta} \mid \tau(\mathbf{y}))
$$

Hence we can compare the approximative and the exact posterior density from

$$
\frac{\pi_{\mathrm{lin}}(\boldsymbol{\theta} \mid \tilde{\boldsymbol{\theta}}(\mathbf{y}))}{\pi(\boldsymbol{\theta} \mid \boldsymbol{y})}=\frac{\pi_{\mathrm{lin}}(\boldsymbol{\theta} \mid \tau(\mathbf{y}))}{\pi(\boldsymbol{\theta} \mid \tau(\mathbf{y}))}
$$

In $\pi_{\text {lin }}(\boldsymbol{\theta} \mid \tau(\mathbf{y}))$ we take $\mathrm{E}_{\boldsymbol{\theta}}[\tau(\mathbf{y})]=\mathbf{J} \boldsymbol{\theta}$, in $\pi(\boldsymbol{\theta} \mid \boldsymbol{\tau}(\mathbf{y}))$ we take $\mathrm{E}_{\boldsymbol{\theta}}[\tau(\boldsymbol{y})]=\boldsymbol{m}_{\boldsymbol{\theta}}$, otherwise the Bayes formulae for computing $\pi_{1 \mathrm{in}}(\boldsymbol{\theta} \mid \tau(\mathbf{y}))$ and $\pi(\boldsymbol{\theta} \mid \tau(\boldsymbol{y}))$ are the same.

## 5. A NOTE ON CONFIDENCE REGIONS FOR $\theta$

We consider confidence regions for $\boldsymbol{\theta}$ which are based on $\tilde{\boldsymbol{\theta}}(\boldsymbol{y})$. We note that they are of restricted importance, since they are influence by the choice of $\boldsymbol{\theta}^{*}$ in (3), resp. by the choice of the points $\boldsymbol{\theta}^{1}, \ldots, \boldsymbol{\theta}^{\boldsymbol{k}}$ in (8) which in fact represents a prior knowledge about $\theta$.

From (16) we obtain that the set

$$
\begin{equation*}
\left\{\boldsymbol{\theta}:\left\|\tilde{\boldsymbol{\theta}}(\mathbf{y})-\mathbf{Q} \boldsymbol{m}_{\boldsymbol{\theta}}\right\|_{\mathbf{V}}^{2}<c \chi_{m}^{2}(\alpha)\right\} \tag{21}
\end{equation*}
$$

is a confidence region for $\boldsymbol{\theta}$ in the case that $c \mathbf{W}$ is known. $\alpha$ is the exact confidence level, and $\chi_{m}^{2}(\alpha)$ is the $\alpha$-quantile of the $\chi^{2}$ distribution with $m$ degrees of freedom.

Example 2. Take the set-up from Example 1. We have

$$
\mathbf{Q} \boldsymbol{m}_{\theta}=\mathrm{E}_{\theta}[\tilde{\theta}(\mathbf{y})]=\cos \delta\left[-\cos \theta \sin \theta^{*}+\sin \theta \cos \theta^{*}\right]+\theta^{*} .
$$

Hence

$$
\left\|\tilde{\theta}(\mathbf{y})-\mathbf{Q} \boldsymbol{m}_{\theta}\right\|_{\mathbf{V}}^{2}=\left[-\left(y_{1}-\cos \theta\right) \sin \theta^{*}+\left(y_{2}-\sin \theta\right) \cos \theta^{*}\right]^{2} .
$$

We see that the confidence region (21) does not depend on $\delta$, hence the standard and the nonstandard linearizations are equivalent as regard to the confidence regions. This is by no way in contradiction to Proposition 2 ; the random variable $\tilde{\theta}(\boldsymbol{y})$ has a small variance, however, this has no importance for confidence reasoning. On the other hand, the obtained confidence region depends very much on $\theta^{*}$.

To understand the situation geometrically, let us write $\tilde{\boldsymbol{\theta}}(\mathbf{y})$ in the form

$$
\tilde{\theta}(\boldsymbol{y})=\mathbf{L y}+\boldsymbol{I}
$$

for some matrix $\mathbf{L}$ and some vector I (This is possible, since $\tilde{\theta}(\mathbf{y})$ is linear in $\mathbf{y}$ ). Further we have

$$
c \mathbf{V}=\operatorname{Var} \tilde{\boldsymbol{\theta}}(\mathbf{y})=c \mathbf{L} \mathbf{W}^{-1} \mathbf{L}^{\mathbf{T}}
$$

hence

$$
\mathbf{P}:=\mathbf{L}^{\mathrm{T}} \mathbf{V}^{-1} \mathbf{L} \mathbf{W}^{-1}
$$

is a $\mathbf{W}$-orthogonal projector. We can verify that

$$
\begin{equation*}
\|\mathbf{P}[\mathbf{y}-\boldsymbol{\eta}(\boldsymbol{\theta})]\|_{\mathbf{w}}^{2}=\left\|\tilde{\theta}(\mathbf{y})-\mathbf{Q} \boldsymbol{m}_{\boldsymbol{\theta}}\right\|_{\mathbf{V}}^{2} \tag{22}
\end{equation*}
$$

Hence the confidence region (21) has the form

$$
\left\{\boldsymbol{\theta}:\|\mathbf{P}[\boldsymbol{y}-\boldsymbol{\eta}(\theta)]\|_{\mathbf{w}}^{\mathbf{w}}<c \chi_{m}^{2}(\alpha)\right\}
$$

This confidence region, although exact, gives poor results (it is too large) if the value $\left\|\mathbf{P}\left[\boldsymbol{\eta}\left(\boldsymbol{\theta}_{\text {true }}\right)-\boldsymbol{\eta}\left(\boldsymbol{\theta}^{*}\right)\right]\right\|_{\mathbf{w}}^{2}$ is large. (We note, that this is zero if model (1) is linear.)

Another consequence of (22) is that $\left\|\tilde{\boldsymbol{\theta}}(\boldsymbol{y})-\mathbf{Q} \boldsymbol{m}_{\boldsymbol{\theta}}\right\|_{\mathbf{V}}^{2}$ and $\|(\mathbf{I}-\mathbf{P})[\boldsymbol{y}-\boldsymbol{\eta}(\boldsymbol{\theta})]\|_{\mathbf{W}}^{2}$ are independent random variables. Hence another confidence region (of the exact confidence level $\alpha$ ) is of the form

$$
\left\{\boldsymbol{\theta}: \frac{(N-m)\left\|\tilde{\boldsymbol{\theta}}(\mathbf{y})-\mathbf{Q} \boldsymbol{m}_{\boldsymbol{\theta}}\right\|_{\mathbf{v}}^{2}}{m\|(\mathbf{I}-\mathbf{P})[\boldsymbol{y}-\boldsymbol{\eta}(\boldsymbol{\theta})]\|_{\mathbf{W}}^{2}}<F_{m, N-m}(\alpha)\right\}
$$

where $\mathrm{F}_{m, N-m}(\alpha)$ is the $\alpha$-quantile of the F-distribution with $m$ and $N-m$ degrees of freedom. The advantage of this region comparing with (21) is that it can be used in the case when $c$ is unknown.

## 6. CONSEQUENCES FOR NONLINEAR EXPERIMENTAL DESIGN

The covariance matrix of $\tilde{\boldsymbol{\theta}}(\mathbf{y})$ (Eq. (16)), and the approximative aposteriori covariance matrix (Eq. (20)) do not depend on the observed vector $\boldsymbol{y}$, and are smaller than the corresponding variances in the standard linearization. Therefore they are adequate to construct optimality criteria for optimum experimental design in nonlinear models.
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