# Lower Bounds for the Minimal Distance in Rational Approximation ${ }^{1}$ 

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## 1. Introduction

We consider a compact Hausdorff space $M$ and denote by $C(M)$ the vector space of real-valued continuous functions on $M$. As a norm in $C(M)$ we introduce the maximum norm

$$
\|g\|_{M}=\max _{x \in M}|g(x)|, \quad g \in C(M) .
$$

Let $U$ and $V$ be finite-dimensional subspaces of $C(M)$ which are spanned by $u_{0}, \ldots, u_{r}$ and $v_{0}, \ldots, v_{s} \in C(M)$. We assume the convex cone

$$
V_{M}{ }^{+}=\{v \in V: v(x)>0 \quad \text { for all } \quad x \in M\}
$$

to be nonempty. For each $f \in C(M)$ we define

$$
\rho_{M}(f)=\inf _{u \in U, v \in V_{M^{+}}}\left\|f-\frac{u}{v}\right\|_{M}
$$

and call this number the minimal distance between $f$ and $W_{M}=\{u / v: u \in U$, $\left.v \in V_{M}{ }^{+}\right\}$.

The rational approximation problem consists of finding $\hat{u} \in U$ and $\hat{v} \in V_{M}{ }^{+}$ such that

$$
\left\|f-\frac{\hat{u}}{\hat{v}}\right\|_{M}=\rho_{M}(f) .
$$

$\hat{u} \mid \hat{v}$ is called a best approximant of $f\left(\mathrm{in} W_{M}\right)$. Each couple $u \in U, v \in V_{M}{ }^{+}$yields a trivial upper bound of $\rho_{M}(f)$. However, for an estimation of the difference between $\|f-(u / v)\|_{M}$ and $\rho_{M}(f)$ or even for an estimation of $\rho_{M}(f)$ itself it is important to know lower bounds of $\rho_{M}(f)$.
In [5] we have developed a principle for the computation of such lower bounds which has been originated by Collatz [2], [3], [4] and has been expanded to nonlinear approximation by Meinardus and Schwedt [10], [11].

In this paper we intend to develop a more general principle which can be handled in a simpler way. For that purpose we consider a nonempty closed

[^0]subset $D$ of $M . D$ is compact and all the functions $g \in C(M)$ can be considered as real-valued continuous functions on $D$ provided with the norm
$$
\|g\|_{D}=\max _{x \in D}|g(x)|
$$

If we define

$$
V_{D}^{+}=\{v \in V: v(x)>0 \quad \text { for all } \quad x \in D\}
$$

we get a nonempty subset of $V$ since $V_{M}{ }^{+}\left(\subseteq V_{D}{ }^{+}\right)$is assumed to be nonempty. Furthermore, for

$$
\rho_{D}(f)=\inf _{u \in U, v \in D_{D^{+}}}\left\|f-\frac{u}{v}\right\|_{D}
$$

we have

$$
\rho_{D}(f) \leqslant \rho_{M}(f)
$$

Our aim is to compute $\rho_{D}(f)$ or at least to find lower bounds for $\rho_{D}(f)$ when $D$ is a certain finite subset of $M$.

The results of this paper, without proofs, have been given in [6].

## 2. Lower Bounds for the Minimal Distance

Notation. By $\boldsymbol{R}^{n}$ we denote the real Euclidean $n$-space and by $\theta_{n}$ the zero vector of $\boldsymbol{R}^{n}$. For $x, y \in \boldsymbol{R}^{n}$ we write $x \geqslant y$ if and only if $x_{i} \geqslant y_{i}$ for $i=1, \ldots, n$. $z^{T}$ denotes the transposed vector $z$. By $|z|$ we mean the vector $\left(\left|z_{1}\right|, \ldots,\left|z_{n}\right|\right)^{T}$, where $z=\left(z_{1}, \ldots, z_{n}\right)^{T}$.

Lemma 2.1. If we assume that for a subset $D=\left\{x_{1}, \ldots, x_{n}\right\}$ of $M$ there exist two vectors $c=\left(c_{1}, \ldots, c_{n}\right)^{T} \neq \theta_{n}$ and $p=\left(p_{1}, \ldots, p_{n}\right)^{T} \geqslant \theta_{n}$ and numbers $\lambda_{1}, \ldots, \lambda_{n} \in \boldsymbol{R}$ such that

$$
\begin{align*}
\sum_{i=1}^{n} u_{j}\left(x_{i}\right) c_{i} & =0, \quad j=0, \ldots, r  \tag{2.1}\\
\sum_{i=1}^{n} f\left(x_{i}\right) v_{k}\left(x_{i}\right) c_{i} & =\sum_{i=1}^{n} \lambda_{i}\left(\left|c_{i}\right|+p_{i}\right) v_{k}\left(x_{i}\right)  \tag{2.2}\\
k & =0, \ldots, s
\end{align*}
$$

then we have

$$
\begin{equation*}
\min _{i=1, \ldots, n} \lambda_{i} \leqslant \rho_{D}(f) \tag{2.3}
\end{equation*}
$$

(The assertion is a slight generalization of Satz 1 in [7] where all the $\lambda_{i}$ 's are assumed to be equal.)

Proof. The case $\lambda_{i} \leqslant 0$ for at least one $i$ is trivial. Hence we assume

$$
\min _{i} \lambda_{i}>0
$$

For a given $u \in U$ and $v \in V_{D}{ }^{+}$the conditions (2.1), (2.2) and $p \geqslant \theta_{n}$ imply

$$
\begin{gathered}
\sum_{i=1}^{n} c_{i}\left(f\left(x_{i}\right)-\frac{u\left(x_{i}\right)}{v\left(x_{i}\right)}\right) v\left(x_{i}\right)=\sum_{i=1}^{n} \lambda_{i}\left(\left|c_{i}\right|+p_{i}\right) v\left(x_{i}\right) \\
\quad \geqslant \sum_{i=1} \lambda_{i}\left|c_{i}\right| v\left(x_{i}\right) \geqslant\left(\min _{i} \lambda_{i}\right) \sum_{i=1}\left|c_{i}\right| v\left(x_{i}\right)
\end{gathered}
$$

and because of

$$
\sum_{i=1}^{n}\left|c_{i}\right| v\left(x_{i}\right)>0
$$

we have

$$
\min \lambda_{i} \leqslant \frac{\sum_{i=1}^{n} c_{i}\left(f\left(x_{i}\right)-\frac{u\left(x_{i}\right)}{v\left(x_{i}\right)}\right) v\left(x_{i}\right)}{\sum_{i=1}^{n}\left|c_{i}\right| v\left(x_{i}\right)} \leqslant\left\|f-\frac{u}{v}\right\|_{D}
$$

Since $u \in U$ and $v \in V_{D}{ }^{+}$are arbitrarily chosen, we can conclude (2.3), which completes the proof.

In the case $\rho_{M}(f)>0$ by Lemma 3.2 of [ 8 ] there exist for each $\lambda \in\left(0, \rho_{M}(f)\right.$ ], $n(\leqslant r+s+3)$ distinct points $x_{1}, \ldots, x_{n}$ and vectors $c \neq \theta_{n}, p \geqslant \theta_{n}$ such that (2.1) and (2.2) hold if we choose

$$
\lambda_{i}=\lambda \quad \text { for } \quad i=1, \ldots, n .
$$

Hence, in principle, $\rho_{M}(f)$ can be estimated from below as best as possible by use of Lemma 2.1. However, Lemma 2.1 is not very convenient for numerical purposes.

In order to find a result which can be handled with less effort we need the following:

Assumption. We require the functions

$$
\begin{equation*}
u_{0}, \ldots, u_{r}, \quad v_{0} \cdot f, \ldots, v_{s} \cdot f \tag{2.4}
\end{equation*}
$$

to be linearly independent on $M$. Under this condition there exist $n=r+s+2$ distinct points $x_{1}, \ldots, x_{n} \in M$ such that the functions $u_{0}, \ldots, u_{r}, v_{0} \cdot f, \ldots, v_{s} \cdot f$ are linearly independent on

$$
D=\left\{x_{1}, \ldots, x_{n}\right\}
$$

This means that the matrix

$$
\begin{equation*}
\tilde{A}=\binom{u_{j}\left(x_{i}\right)}{v_{k}\left(x_{i}\right) f\left(x_{i}\right)}^{2} \tag{2.5}
\end{equation*}
$$

is nonsingular. If we define the matrix $B$ by

$$
\begin{equation*}
B=\binom{O}{v_{k}\left(x_{i}\right)} \tag{2.6}
\end{equation*}
$$

where $O$ is a zero matrix consisting of $r+1$ rows and $r+s+2$ columns we can formulate

Lemma 2.2. We assume $D=\left\{x_{1}, \ldots, x_{n}\right) \subseteq M, n=r+s+2$, to be such that the matrix (2.5) is nonsingular. Then for each vector $y=\left(y_{1}, \ldots, y_{n}\right)^{T} \geqslant \theta_{n}$, $y \neq \theta_{n}$, there exists exactly one vector $c=\left(c_{1}, \ldots, c_{n}\right)^{T} \neq \theta_{n}$ such that

$$
\begin{equation*}
\tilde{A c}=B y \tag{2.7}
\end{equation*}
$$

and we have

$$
\begin{equation*}
q(y)=\min _{c_{i} \neq 0} \frac{y_{i}}{\left|c_{i}\right|} \leqslant \rho_{D}(f) . \tag{2.8}
\end{equation*}
$$

Remark. If we choose $y>\theta_{n}$ (that is $y_{i}>0$ for $i=1, \ldots, n$ ), we have $q(y)>0$. Hence in this case we always get a positive lower bound of $\rho_{M}(f)$ by solving the linear system (2.7) and computing $q(y)>0$. Furthermore, the assumption (2.4) yields $\rho_{M}(f)>0$.

Proof. For each $y \in \boldsymbol{R}^{n}, y \geqslant \theta_{n}, y \neq \theta_{n}$ we have

$$
B y \neq \theta_{n}
$$

Otherwise, for each $v \in V_{D}{ }^{+}$we would have

$$
\sum_{i=1}^{n} v\left(x_{i}\right) y_{i}=0
$$

which is impossible. Since $\tilde{A}$ is nonsingular for each such $y \in \boldsymbol{R}^{n}$ there is exactly one solution $c \neq \theta_{n}$ of the linear system (2.7). We put $I=\left\{i: c_{i}=0\right.$ and $\left.y_{i}=0\right\}$ and define for each $i \notin I$

$$
\lambda_{i}=\left\{\begin{array}{lll}
\frac{y_{i}}{\left|c_{i}\right|} & \text { if } & c_{i} \neq 0 \\
\frac{y_{i}}{p_{i}} & \text { if } & c_{i}=0
\end{array}\right.
$$

[^1]where $p_{i}>0$ is at our disposal. If we put $p_{i}=0$ for each $i$ such that $c_{i} \neq 0$, the system (2.7) can be written as
\[

$$
\begin{aligned}
& \sum_{i \notin I} u_{j}\left(x_{i}\right) c_{i}=0, \quad j=0, \ldots, r \\
& \sum_{i \neq I} f\left(x_{i}\right) v_{k}\left(x_{i}\right) c_{i}=\sum_{i \notin I} \lambda_{i}\left(\left|c_{i}\right|+p_{i}\right) v_{k}\left(x_{i}\right), \quad k=0, \ldots, s .
\end{aligned}
$$
\]

Hence by Lemma 2.1 we conclude

$$
\min _{i \notin I} \lambda_{i} \leqslant \rho_{D}(f) .
$$

If for each $i \notin I$ such that $c_{l}=0$ we choose $p_{l}>0$ sufficiently small, we can achieve

$$
q(y)=\min _{c_{i} \neq 0} \frac{y_{i}}{\left|c_{l}\right|}=\min _{i \notin I} \lambda_{i}
$$

which completes the proof.
Theorem 2.1. Under the assumption of Lemma 2.2 for the set $D \subseteq M$ we have

$$
\rho_{D}(f)=\max _{y \in K} q(y),
$$

where $q(y)$ is defined by (2.8) and

$$
\begin{equation*}
K=\left\{y \in \boldsymbol{R}^{n}: y \geqslant \theta_{n}, \quad y \neq \theta_{n}\right\} . \tag{2.9}
\end{equation*}
$$

Proof. We have to show that there exists $\hat{y} \in K$ such that

$$
q(\hat{y})=\rho_{D}(f) .
$$

Then the assertion follows by Lemma 2.2. By Satz 2 in [7] there exist $\hat{c} \in \boldsymbol{R}^{n}$, $\hat{c} \neq \theta_{n}$ and $p \in R^{n}, p \geqslant \theta_{n}$ so that

$$
\tilde{A} \hat{c}=\rho_{D}(f) \cdot B(|\hat{c}|+p)
$$

where $\tilde{A}$ and $B$ are given by (2.5) and (2.6). If we put $\hat{y}=\rho_{D}(f)(|\hat{c}|+p)$, then $\hat{y} \in K$, and by Lemma 2.2

$$
q(\hat{y})=\min _{\hat{c}_{1} \neq 0} \frac{\hat{y}_{i}}{\left|\hat{c}_{i}\right|} \leqslant \rho_{D}(f) .
$$

On the other hand, for each $i$ such that $\hat{c}_{i} \neq 0$ we have

$$
\frac{\hat{y}_{i}}{\left|\hat{c}_{i}\right|}=\rho_{D}(f) \frac{\left|\hat{c}_{i}\right|+p_{i}}{\left|\hat{c}_{i}\right|} \geqslant \rho_{D}(f),
$$

whence $\rho_{D}(f) \leqslant q(\hat{y})$, which completes the proof.

In the setting of Theorem 2.1 the computation of $\rho_{D}(f)$ leads to the following nonlinear optimization problem (which is solvable): Under the conditions

$$
|A y| \leqslant \frac{1}{\lambda} y, \quad y \geqslant \theta_{n}, \quad y \neq \theta_{n}
$$

$\lambda$ is to be maximized $\left(A=\tilde{A}^{-1} B\right)$.

## 3. A Nonlinear Eigenvalue Problem

Let $D=\left\{x_{1}, \ldots, x_{n}\right\}, n=r+s+2$, be a subset of $M$ such that the matrix (2.5) is nonsingular. We consider the following problem: Find a number $\lambda>0$ such that the system

$$
\begin{gather*}
\sum_{i=1}^{n} u_{j}\left(x_{i}\right) c_{i}=0, \quad j=0, \ldots, r  \tag{3.1}\\
\sum_{i=1}^{n} v_{k}\left(x_{i}\right) f\left(x_{i}\right) c_{i}=\lambda \sum_{i=1}^{n}\left|c_{i}\right| v_{k}\left(x_{i}\right) \\
k=0, \ldots, s
\end{gather*}
$$

has a solution $c=\left(c_{i}, \ldots, c_{n}\right)^{T} \neq \theta_{n}$ For each such number $\lambda$ we have, by Lemma 2.1,

$$
\lambda \leqslant \rho_{\mathrm{D}}(f)
$$

If there is a best approximant of $f$ in

$$
W_{D}=\left\{\frac{u}{v} ; u \in U, v \in V_{D}^{+}\right\} \subseteq C(D)
$$

we know, [7], that for $\lambda=\rho_{D}(f)$ there is a nontrivial solution $c$ of (3.1). Furthermore, we know by Satz 2.1 in [8] that there is a subset $D$ of $M$ such that for $\lambda=\rho_{M}(f)$ the system (3.1) admits a nontrivial solution $c$ if the approximation problem in $C(M)$ is solvable. By the substitution

$$
\begin{equation*}
y=\lambda|c|, \quad \mu=\frac{1}{\lambda} \tag{3.2}
\end{equation*}
$$

the above problem turns out to be equivalent to the following nonlinear eigenvalue problem: Find a number $\mu>0$ such that there is a solution $y \in K$ of

$$
\begin{equation*}
|A y|=\mu y \tag{3.3}
\end{equation*}
$$

where $A=\tilde{A}^{-1} B, \tilde{A}$ is given by (2.5), $B$ by (2.6) and $K$ by (2.9).

Theorem 3.1. There is $a y \in K$ and $a \mu>0$ such that (3.3) holds, i.e., such that

$$
c=A y \quad \text { and } \quad \lambda=\frac{1}{\mu}
$$

solve problem (3.1).
Proof. (According to [9].) We define a mapping $P: \boldsymbol{R}^{\boldsymbol{n}} \rightarrow \boldsymbol{R}^{\boldsymbol{n}}$ by

$$
\begin{equation*}
P(y)=|A y|, \quad y \in R^{n} \tag{3.4}
\end{equation*}
$$

In the proof of Lemma 2.2 we have shown

$$
B y \neq \theta_{n} \quad \text { for each } \quad y \in K .
$$

Hence $P(y)=|A y|=\left|\tilde{A}^{-1} B y\right| \geqslant \theta_{n}$ and $\neq \theta_{n}$ for each $y \in K$; that is, $P(K) \subseteq K$. Evidently the subset

$$
S=\left\{y \in K:\|y\|_{1}:=\sum_{i=1}^{n} y_{i}=1\right\}
$$

of $K$ is convex and compact and we have

$$
\|P(y)\|_{1}:=\sum_{i=1}^{n} P(y)_{i}>0
$$

for all $y \in S$. Hence the operator

$$
\tilde{P}(y)=\frac{P(y)}{\|P(y)\|_{1}}
$$

is defined on $S$, continuous and maps $S$ into itself. By Brouwer's fixed-point theorem there is a $y^{*} \in S$ such that $\widetilde{P}\left(y^{*}\right)=y^{*}$ or

$$
P\left(y^{*}\right)=\left|A y^{*}\right|=\mu^{*} y^{*}
$$

where $\mu^{*}=\left\|P\left(y^{*}\right)\right\|_{1}>0$. This completes the proof.
For numerical purposes it would be very helpful if the operator $P$ defined by (3.4) were monotone on $K \cup\left\{\theta_{n}\right\}$; that is,

$$
\theta_{n} \leqslant y \leqslant z \quad \text { implies } \quad P(y) \leqslant P(z) .
$$

If we then start with an arbitrary $y^{0} \in K$ and define a sequence $y^{k} \in K$ by

$$
y^{k+1}=P\left(y^{k}\right), \quad k=0,1,2, \ldots
$$

it turns out that for the numbers

$$
q_{k}=\min _{y_{i}^{k+1} \neq 0} \frac{y_{i}^{k}}{y_{i}^{k+1}} \quad \text { and } \quad \hat{q}_{k}=\max _{y_{i}^{k+1} \neq 0} \frac{y_{i}^{k}}{y_{i}^{k+1}}
$$

we have

$$
q_{0} \leqslant q_{1} \leqslant \ldots \leqslant q_{k} \leqslant \rho_{D}(f) \leqslant \hat{q}_{k} \leqslant \ldots \leqslant \hat{q}_{1} \leqslant \hat{q}_{0}
$$

and

$$
\rho_{D}(f)=\lim _{k \rightarrow \infty} q_{k}=\lim _{k \rightarrow \infty} \hat{q}_{k}
$$

if $P$ is strictly monotone in the sense of Bohl [l].
If we define the matrix $|A|$ by taking the absolute values of the elements of $A$ as elements of $|A|$, then $P$ is obviously monotone on $K \cup\left\{\theta_{n}\right\}$ if we have

$$
\begin{equation*}
|A y|=|A| y \quad \text { for each } \quad y \in K \tag{3.5}
\end{equation*}
$$

Sufficient for (3.5) is that in each row of $A$ all the elements which are unequal to zero have the same sign. But as E. Bohl pointed out this is also necessary for $P$ to be monotone on $K \cup\left\{\theta_{n}\right\}$. Bohl gave a simple (unpublished) proof for this fact, namely: Assume that for some index $i$ and two indices $j$ and $k$ with $j \neq k$ we have $A_{i j} \neq 0, A_{i k} \neq 0$ and $\operatorname{sgn} A_{i j}=-\operatorname{sgn} A_{i k}$. Then we define two vectors $y$ and $z \in \boldsymbol{R}^{n}$ by

$$
\begin{gathered}
y_{l}=\left\{\begin{array}{cc}
0 & \text { for } l \neq j \text { and } k, \\
\lambda & \text { for } l=k, \\
-\frac{A_{i k}}{A_{i j}} & \text { for } l=j,
\end{array}\right\} \text { and } z_{l}=\left\{\begin{array}{ccl}
0 & \text { for } l \neq j \text { and } k, \\
1 & \text { for } l=k, \\
-\frac{A_{i k}}{A_{i j}} & \text { for } l=j,
\end{array}\right\} \\
l=1, \ldots n .
\end{gathered}
$$

If we choose $0 \leqslant \lambda<1$, then $y, z \in K$ and

$$
\theta_{n} \leqslant y \leqslant z
$$

However,

$$
\begin{aligned}
& \left|(A y)_{i}\right|=\left|-A_{i k}+\lambda A_{i k}\right|=\left|A_{i k}\right|(1-\lambda) \\
& \left|(A z)_{i}\right|=\left|-A_{i k}+A_{i k}\right|=0
\end{aligned}
$$

and $\left|(A y)_{i}\right| \leqslant\left|(A z)_{i}\right|=0$ would imply $A_{i k}=0$, a contradiction. The result is that the operator $P$ defined by (3.4) is monotone if and only if the problem (3.3) is a linear eigenvalue problem of the form (3.5).

In general, $P$ is not monotone as the following example shows: $M=[a, b]$, $r=0, s=1, u_{0}=v_{0} \equiv 1, v_{1}(x)=x$ and $f \in C(M)$ chosen such that for three different points $x_{i} \in[a, b], i=1,2,3$, we have $f\left(x_{1}\right) \neq 0, f\left(x_{3}\right) \neq 0, f\left(x_{2}\right)=0$. Then the matrix $\tilde{A}$ given by (2.5) is nonsingular in this case; in fact, the determinant of $\tilde{A}$ is given by the formula

$$
\operatorname{det}(\widetilde{A})=\left(x_{1}-x_{3}\right) f\left(x_{1}\right) f\left(x_{3}\right)
$$

and hence is not equal to zero. $A=\tilde{A}^{-1} B$, with $B$ of (2.6), is given by

$$
A=\left(\begin{array}{ccc}
\frac{1}{f\left(x_{1}\right)} & \frac{x_{2}-x_{3}}{\left(x_{1}-x_{3}\right) f\left(x_{1}\right)} & 0 \\
\frac{1}{f\left(x_{1}\right)} & \frac{\left(x_{2}-x_{1}\right) f\left(x_{1}\right)+\left(x_{3}-x_{2}\right) f\left(x_{3}\right)}{\left(x_{1}-x_{3}\right) f\left(x_{1}\right) f\left(x_{3}\right)} & \frac{1}{f\left(x_{3}\right)} \\
0 & \frac{x_{1}-x_{2}}{\left(x_{1}-x_{3}\right) f\left(x_{3}\right)} & \frac{1}{f\left(x_{3}\right)}
\end{array}\right)
$$

Therefore $P$ defined by (3.4) is monotone on $K \cup\left\{\theta_{n}\right\}$ if and only if

$$
\operatorname{sgn} f\left(x_{1}\right)=\operatorname{sgn} f\left(x_{3}\right),
$$

and

$$
\operatorname{sgn}\left(x_{1}-x_{2}\right)=\operatorname{sgn}\left(x_{1}-x_{3}\right)=\operatorname{sgn}\left(x_{2}-x_{3}\right) .
$$

Final remark: In [6] we have shown how the nonlinear eigenvalue problem (3.3) is related to the linear eigenvalue problem which has been investigated by Werner [12] in the case of the classical rational approximation problem.

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[^1]:    ${ }^{2} j=0, \ldots, r$ and $k=0, \ldots, s$ denote row indices and $i=1, \ldots, r+s+2$ denotes column indices.

