# A PRIORI BOUNDS FOR THE RICCATI EQUATION 

R. S. BUCY<br>University of Southern California

## 1. Introduction and general results

Since the appearance of [6] and [7], the theory of linear filtering has experienced a renaissance. This theory, although evidently well known to statisticians in terms of "least squares estimates," has found many applications in the early sixties, largely because of the realization and synthesis methods provided in [6] and [7]. For an indication of some of the aerospace applications to guidance of spacecraft, the interested reader may find detailed information in [3]. Although the theory of linear filtering has changed little from that given in [7] for the continuous time problem, the practical realization of the so-called "correlated noise problem" as treated mathematically in [6], has recently found a solution in [4]. The full solution of this discrete time filtering problem and its meaning is described in detail in [5]. For readers desirous of a survey of recent results in linear and nonlinear filtering, it is available in [5], while more detailed information can be found in [3].

In this paper, our interest will center on the discrete matrix Riccati equation with emphasis on the study of the asymptotic behavior of its covariance matrix solution. A major tool in this study will be the Duffin parallel resistance of two nonnegative definite matrices $A$ and $B$ denoted by $A: B$. This operation is described in detail in [1] and provides for us a link between the Riccati equation and the classical continued fraction theory described in [9] and [10].

We have undertaken to study the discrete Riccati equation from the point of view of continued fractions because this technique provides considerable generality in that the nonsingular theory becomes a rather special case (see [3], Chapter 5) and much deeper results are obtained for singular problems; also the methods are striking generalizations of classical continued fraction methods.

We will be concerned with the cone $C$ of $d \times d$ real entry symmetric nonnegative definite matrices. The cone $C$ induces a natural partial ordering as for $A \in C$ and $B \in C, A \geqq B$ when and only when $A-B \in C$. The object of study will be the map of $C$;

This research was supported in part by the United States Air Force, Office of Aerospace Research, Applied Mathematics Division, under Grant AF-AF OSR-1244-67.

$$
\begin{equation*}
\tau_{n}(A)=\phi_{n} H_{n-1}^{\prime}\left(H_{n-1} A H_{n-1}^{\prime}: R_{n-1}\right) H_{n-1} \phi_{n}^{\prime}+C_{n}, \tag{1.1}
\end{equation*}
$$

where with $M_{\ell, m}(R)$ the set of $\ell \times m$ real entry matrices.
The matrices $A$ and $\phi_{n}$ belong to $M_{d, d}(R), H_{n-1} \in M_{s, d}(R), R_{n-1} \in M_{s, s}(R)$, and $R_{n-1}$ is positive definite. The matrix $C_{n} \in M_{d, d}(R)$ with $s \geqq d$ and $H_{n-1}^{\prime} H_{n-1}=I_{d}$, the $d \times d$ identity. We will restrict $C_{n}$ and $A$ to be members of the cone $C$ and will center our attention on iterates of the mapping $\tau$. We will consider the case where $s \geqq d$ and provide motivation for (1.1) in terms of the filtering problem in a later section. We recall that $(A: B)=A(A+B)^{*} B$ (see [1]), where ${ }^{\text {\# }}$ denotes the Moore-Penrose pseudoinverse determined by the axioms:
(a) $A A^{\#} A=A$,
(b) $A^{\#} A A^{\#}=A^{\#}$,
(c) $\left(A A^{\#}\right)^{\prime}=A A^{\#}$,
(d) $\left(A^{\#} A\right)^{\prime}=A^{\#} A$,
with ' denoting transpose.
The following series of lemmas will prove useful in the sequel.
Lemma 1.1. For $A_{i} \in M_{\ell, k}(R)$,

$$
\begin{equation*}
A_{i}\left[\sum_{j=1}^{n} A_{j}^{\prime} A_{j}\right]\left[\sum_{j=1}^{n} A_{j}^{\prime} A_{j}\right]^{\#}=A_{i} \tag{1.2}
\end{equation*}
$$

for $1 \leqq i \leqq n$.
Proof. See [8].
Lemma 1.2. For $A$ and $B \in C$ and $A \geqq B, \tau_{n}(A) \geqq \tau_{n}(B)$, and in particular $\tau_{n}(C) \cong C$.

Proof. Denoting by $\|x\|_{A}^{2}$ the quadratic form induced by $A \in C$, it is quite easily seen that

$$
\begin{equation*}
\|\mathbf{x}\|_{\tau_{n}(A)}^{2}=\min _{\mathbf{r} \in R^{s}}\left\{\left\|\phi_{n}^{\prime} \mathbf{x}+H_{n-1}^{\prime} \mathbf{r}\right\|_{A}^{2}+\|\mathbf{r}\|_{R_{n-1}}^{2}+\|\mathbf{x}\|_{C_{n}}^{2}\right\} \tag{1.3}
\end{equation*}
$$

so that

$$
\begin{equation*}
\tau_{n}(A)=\phi_{n}\left(A-A H_{n-1}^{\prime}\left(H_{n-1} A H_{n-1}^{\prime}+R_{n-1}\right)^{\#} H_{n-1} A\right) \phi_{n}^{\prime}+C_{n} \tag{1.4}
\end{equation*}
$$

In particular, if $A \geqq B$, then

$$
\begin{align*}
& \|\mathbf{x}\|_{\tau_{n}(A)}^{2}  \tag{1.5}\\
& \quad=\min _{\mathbf{r} \in R^{s}}\left\{\left\|\phi_{n}^{\prime} \mathbf{x}+H_{n-1}^{\prime} \mathbf{r}\right\|_{B}^{2}+\|\mathbf{r}\|_{R_{n-1}}^{2}+\|\mathbf{x}\|_{C_{n}}^{2}+\left\|\phi_{n}^{\prime} \mathbf{x}+H_{m}^{\prime} \mathbf{r}\right\|_{A-B}^{2}\right\}
\end{align*}
$$

the result follows. Of course, if $S_{n-1}(A)=A-A H_{n-1}^{\prime}\left(H_{n-1} A H_{n-1}^{\prime}+\right.$ $\left.R_{n-1}\right)^{\#} H_{n-1} A$, then

$$
\begin{align*}
& H_{n-1} S_{n-1}(A) H_{n-1}^{\prime}  \tag{1.6}\\
& \quad=H_{n-1} A H_{n-1}^{\prime}\left(H_{n-1} A H_{n-1}^{\prime}+R_{n-1}\right)^{\#} R_{n-1}=\left(H_{n-1} A H_{n-1}\right): R_{n-1}
\end{align*}
$$

in view of Lemma 1.1, or

$$
\begin{equation*}
S_{n-1}(A)=H_{n-1}^{\prime}\left(R_{n-1}-R_{n-1}\left(H_{n-1} A H_{n-1}^{\prime}+R_{n-1}\right)^{\#} R_{n-1}\right) H_{n-1} \tag{1.7}
\end{equation*}
$$

since by our assumption, $H_{n-1}^{\prime} H_{n-1}=I_{d}$.

We will be concerned with the iterated composition mapping $\tau_{1}\left(\tau_{2}(\cdots\right.$ $\left.\left(\tau_{n}(A)\right) \cdots\right)$. It is clear that $\tau_{1} \tau_{2} \cdots \tau_{n}(0)$ is monotone increasing in $C$, while $\tau_{1} \cdots \tau_{n}(\infty)$ is monotone decreasing in $C$. We remark that Riesz has shown bounded monotone sequences in $C$ converge, that is, if $A_{n} \geqq A_{n-1}$ and $A_{n} \leqq \alpha I_{d}$ then $\lim _{n \rightarrow \infty} A_{n}$ exists.

For a direct generalization of the scalar continued fraction theory see Wall [10].

Theorem 1.1. Let $t_{i}(A)=D_{1}+A$ if $i=1$, and

$$
\begin{equation*}
t_{i}(A)=-\phi_{i-1} H_{i-2}^{\prime} R_{i-2}\left(H_{i-2}\left(A+D_{i}\right) H_{i-2}^{\prime}+R_{i-2}\right)^{\#} R_{i-2} H_{i-2} \phi_{i-1}^{\prime} \tag{1.8}
\end{equation*}
$$

otherwise then

$$
\begin{align*}
& \tau_{1} \tau_{2} \cdots \tau_{n}(A)  \tag{1.9}\\
& \quad=t_{1} \cdots t_{n}\left(-\phi_{n} H_{n-1}^{\prime} R_{n-1}\left(H_{n-1} A H_{n-1}^{\prime}+R_{n-1}\right)^{\#} R_{n-1} H_{n-1} \phi_{n}^{\prime}\right)
\end{align*}
$$

where $D_{i}=C_{i}+\phi_{i} H_{i-1}^{\prime} R_{i-1} H_{i-1} \phi_{i}^{\prime}$.
Proof. An induction proof will be used to establish (1.9). For $n=1$, it follows that,

$$
\begin{align*}
\tau_{1}(A) & =D_{1}-\phi_{1} H_{0}^{\prime} R_{0}\left(H_{0} A H_{0}^{\prime}+R_{0}\right)^{\#} R_{0} H_{0} \phi_{1}^{\prime}  \tag{1.10}\\
& =t_{1}\left(-\phi_{1} H_{0}^{\prime} R_{0}\left(H_{0} A H_{0}^{\prime}+R_{0}\right)^{\#} R_{0} H_{0} \phi_{1}^{\prime}\right)
\end{align*}
$$

Now assume (1.9) is true for all $n \leqq k$, then

$$
\begin{align*}
& \tau_{1} \cdots \tau_{k+1}(A)  \tag{1.11}\\
& =\tau_{1} \cdots \tau_{k}\left(\tau_{k+1}(A)\right) \\
& =t_{1} \cdots t_{k}\left(-\phi_{k} H_{k-1}^{\prime} R_{k-1}\left(H_{k-1} \tau_{k+1}(A) H_{k-1}^{\prime}+R_{k-1}\right)^{\#} R_{k-1} H_{k-1} \phi_{k}^{\prime}\right)
\end{align*}
$$

in view of the induction hypothesis. Now by the definition of $t_{k+1},(1.11)$ is equal to $t_{1} \cdots t_{k+1}\left(\tau_{k+1}(A)-D_{k+1}\right)$; however,

$$
\begin{equation*}
\tau_{k+1}(A)-D_{k+1}=-\phi_{k+1} H_{k}^{\prime} R_{k}\left(H_{k} A H_{k}^{\prime}+R_{k}\right)^{*} R_{k} H_{k} \phi_{k+1}^{\prime} \tag{1.12}
\end{equation*}
$$

in view of the proof of Lemma 1.2, so that the assertion follows.
Corollary 1.1. For $A \in C$,

$$
\begin{align*}
\tau_{1} \cdots \tau_{n}(0) & =t_{1} \cdots t_{n}\left(-\phi_{n} H_{n-1}^{\prime} R_{n-1} H_{n-1} \phi_{n}^{\prime}\right)  \tag{1.13}\\
& \leqq \tau_{1} \cdots \tau_{n}(A) \leqq t_{1} \cdots t_{n}(0)=\tau_{1} \cdots \tau_{n}(\infty)
\end{align*}
$$

Proof. The proof is immediate.
Corollary 1.2. Let $t_{i}^{*}(A)=D_{i}+$ A. The limits $V_{i}=\lim _{n \rightarrow \infty} t_{i}^{*} t_{i+1} \cdots t_{n}(0)$ and $L_{i}=\lim _{n \rightarrow \infty} t_{i}^{*} t_{i+1} \cdots t_{n}\left(-\phi_{n} H_{n-1}^{\prime} R_{n-1} H_{n-1} \phi_{n}^{\prime}\right)$ exist and $L_{i} \leqq V_{i}$.

Proof. The map $t_{i}^{*} t_{i+1} \cdots t_{n}(\infty)=\lim _{A \dagger \infty} \tau_{i} \cdots \tau_{n}(A)$ is monotone nonincreasing in $n$ and bounded above by $D_{1}$ and below by 0 and hence converges. A similar argument demonstrates the existence of the other limit.

Remark 1.1. Notice that $V_{i}$ and $L_{i}$ are equilibrium solutions in that from the relation

$$
\begin{equation*}
\tau_{i} \tau_{i+1} \cdots \tau_{n}(A)=\tau_{i}\left(\tau_{i+1} \cdots \tau_{n}(A)\right) \tag{1.14}
\end{equation*}
$$

and as $\tau_{i}$ preserves order in $C$, it follows that $V_{i}=\tau_{i}\left(V_{i+1}\right)$ and $L_{i}=\tau_{i}\left(L_{i+1}\right)$ when $L_{i} \in \dot{C}, \dot{C}$ the interior of $C$. Further, when $\tau_{i}$ is autonomous, $L_{i}$ and $V_{i}$ are constant.

Corollary 1.3. Suppose $\left[\Pi_{j=i}^{n-1} \rho_{j, j+1}^{n}\right]\left\|D_{n}-C_{n}\right\|$ tends to zero as $n \rightarrow \infty$, where

$$
\begin{equation*}
\rho_{i, i+1}^{n}=\left\|\phi_{i} H_{i-1}^{\prime} R_{i-1}\left[\left(H_{i-1} L_{i+1, n} H_{i-1}+R_{i-1}\right)^{\#}\right]^{2} R_{i-1} H_{i-1} \phi_{i}^{\prime}\right\| \tag{1.15}
\end{equation*}
$$

and $L_{i, n}=\tau_{i} \cdots \tau_{n}(0)$, then
(1.16) $L_{i}=\lim _{n \rightarrow \infty} \tau_{i} \cdots \tau_{n}(0)=\lim _{n \rightarrow \infty} \tau_{i} \cdots \tau_{n}(A)=\lim _{n \rightarrow \infty} \tau_{i} \cdots \tau_{n}(\infty)=V_{i}$
for all $A \in C$.
Proof. Let $E_{i, n}=\tau_{i} \cdots \tau_{n}(\infty)-\tau_{i} \cdots \tau_{n}(0)$; then with $E_{n, n}=D_{n}-C_{n}$ and with $L_{i, n}=\tau_{i} \cdots \tau_{n}(0)$

$$
\begin{align*}
& E_{i, n}=\phi_{i} H_{i-1}^{\prime}\left\{\left(H_{i-1}\left[E_{i+1, n}+L_{i+1, n}\right] H_{i-1}^{\prime}: R_{i-1}\right)\right.  \tag{1.17}\\
&\left.-\left(H_{i-1} L_{i+1, n} H_{i-1}^{\prime}: R_{i-1}\right)\right\} H_{i-1} \phi_{i}^{\prime}
\end{align*}
$$

or in view of Lemma 27 in [1]

$$
\begin{equation*}
E_{i, n}=\phi_{i} H_{i-1}^{\prime} R_{i-1} C_{i, n}^{\#}\left(E_{i+1, n}: C_{i, n}\right) C_{i, n}^{\#} R_{i-1} H_{i-1} \phi_{i}^{\prime}, \tag{1.18}
\end{equation*}
$$

where $C_{i, n}=H_{i-1} L_{i+1, n} H_{i-1}^{\prime}+R_{i-1}$. Hence, if $\varepsilon_{i, n}=\left\|E_{i, n}\right\|$, then $\varepsilon_{i, n} \leqq$ $\rho_{i, i+1}^{n} \varepsilon_{i+1, n}$, where

$$
\begin{equation*}
\rho_{i, i+1}^{n}=\left\|\phi_{i} H_{i-1}^{\prime} R_{i-1}\left[\left(H_{i-1} L_{i+1, n} H_{i-1}^{\prime}+R_{i-1}\right)^{\#}\right]^{2} R_{i-1} H_{i-1} \phi_{i}^{\prime}\right\| \tag{1.19}
\end{equation*}
$$

so that

$$
\begin{equation*}
\varepsilon_{i, n} \leqq\left[\prod_{j=i}^{n-1} \rho_{j, j+1}^{n}\right]\left\|D_{n}-C_{n}\right\| \tag{1.20}
\end{equation*}
$$

by Corollary 1.1.
Example 1.1. Consider the scalar Riccati equation

$$
\begin{equation*}
p_{n}=\alpha_{n}^{2} p_{n-1}-\frac{\alpha_{n}^{2} p_{n-1}^{2}}{r_{n-1}+p_{n-1}}+q_{n} \tag{1.21}
\end{equation*}
$$

where $q_{n} \geqq 0$ and $r_{n} \geqq 0$. The associated continued fraction is

$$
\begin{equation*}
b_{0}+\frac{a_{1} \mid}{\mid b_{1}}+\frac{a_{2} \mid}{\mid b_{2}} \cdots \tag{1.22}
\end{equation*}
$$

where

$$
\begin{aligned}
b_{0} & =q_{0}+\alpha_{0}^{2} r_{-1} \\
b_{i} & =q_{i}+\alpha_{i}^{2} r_{i-1}+r_{i-2} \\
a_{i} & =\alpha_{i-1}^{2} r_{i-2}^{2}
\end{aligned}
$$

In view of Corollary 1.3, it suffices that

$$
\begin{equation*}
\left[\prod_{j=1}^{n-1} \frac{\alpha_{j}^{2} r_{j-1}^{2}}{\left(q_{j+1}+r_{j-1}\right)^{2}}\right] \alpha_{n}^{2} r_{n-1} \tag{1.24}
\end{equation*}
$$

tends to zero as $n \rightarrow \infty$. Notice this result is far stronger than those of [2]. In the autonomous case the associated continued fraction is

$$
\begin{equation*}
y=\beta_{0}^{*}+\frac{\alpha^{*} \mid}{\mid \beta^{*}}+\frac{\alpha^{*} \mid}{\mid \beta^{*}}+\cdots, \tag{1.25}
\end{equation*}
$$

where $\beta_{0}^{*}=\beta^{*}-r=q+\alpha^{2} r$ and $\alpha^{*}=-\alpha^{2} r^{2}$ and it is well known that the continued fraction converges to $y=x-r$ with $x$ the positive root of $x^{2}-\beta^{*} x-\alpha^{*}=0$. For detailed information see [10], problem 1.4, page 23. Explicitly, $y$ is given by

$$
\begin{equation*}
y=\frac{1}{2}\left(q+\left(\alpha^{2}+1\right) r+\left[\left(q+\left(\alpha^{2}+1\right) r\right)^{2}-4 \alpha^{2} r\right]^{1 / 2}\right)-r \tag{1.26}
\end{equation*}
$$

Further, if $\alpha=e^{\Delta f}, r=r^{*} / \Delta$, and $q=q^{*} \Delta$, it is easily verified that as $\Delta \rightarrow 0$

$$
\begin{equation*}
y(\Delta) \downarrow f r^{*}+\left(f^{2}\left(r^{*}\right)^{2}+r^{*} q^{*}\right)^{1 / 2} \tag{1.27}
\end{equation*}
$$

the equilibrium solution of the analogous continuous time Riccati equation (see Example 1, p. 98, [7]).

Corollary 1.4. For $n \geqq 1, C_{1} \leqq \tau_{1} \cdots \tau_{n}(A) \leqq D_{1}$.
Theorem 1.2. If $S_{\alpha}^{+}=\left\{A \in C \mid \tau_{n}(A) \geqq A\right.$ for all $\left.n \geqq \alpha\right\}$ and $S_{\alpha}^{-}=$ $\left\{A \in C \mid \tau_{n}(A) \leqq A\right.$ for all $\left.n \geqq \alpha\right\}$, then $\tau_{\alpha} \cdots \tau_{n}(A)$ is monotone nondecreasing for $A \in S^{+}$and monotone nonincreasing for $A \in S^{-}$. Further, if there exist real $\gamma_{\alpha}$ and $\beta_{\alpha}$ such that $D_{n} \leqq \gamma_{\alpha} I, C_{n} \geqq \beta_{\alpha} I$, for all $n \geqq \alpha$, then

$$
\begin{align*}
& S_{\alpha}^{+} \supseteqq\left\{A \in C \mid A \leqq \beta_{\alpha} I\right\},  \tag{1.28}\\
& S_{\alpha}^{-} \supseteqq\left\{A \in C \mid A \geqq \gamma_{\alpha} I\right\} .
\end{align*}
$$

Proof. Since $\tau_{n}$ preserves the ordering of $C$, if $\tau_{n}(A) \geqq A$, then $\tau_{n-1} \tau_{n}(A) \geqq$ $\tau_{n-1}(A)$, and so forth, so that $\tau_{\alpha} \cdots \tau_{n}(A) \geqq \tau_{\alpha} \cdots \tau_{n-1}(A)$ for $A \in S_{\alpha}^{+}$. The other assertions are obvious.

Remark 1.2. In view of the double argument convexity of $(A: B), S_{\alpha}^{+}$is convex (see [1], Theorem 24).

We will now specialize further and assume that $\phi_{i}, C_{i}$, and $R_{i}$ are invertible for all $i$; in view of Corollary 1.4, it suffices to consider $\tau_{n}(A)$ for $A \in C$ the interior of $C$.

Now

$$
\begin{align*}
\tau_{n}(A) & =\phi_{n}\left(A^{-1}+O_{n}^{-1}\right)^{-1} \phi_{n}^{\prime}+C_{n}  \tag{1.29}\\
& =\phi_{n}\left(A: O_{n}\right) \phi_{n}^{\prime}+C_{n}
\end{align*}
$$

with $O_{n}=\left(H_{n-1}^{\prime} R_{n-1}^{-1} H_{n-1}\right)^{-1}$. This relation follows since

$$
\begin{align*}
& \left(H_{n-1} A H_{n-1}^{\prime}: R_{n-1}\right)  \tag{1.30}\\
& \quad=H_{n-1} A H_{n-1}^{\prime}\left(H_{n-1} A H_{n-1}^{\prime}+R_{n-1}\right)^{-1} R_{n-1} \\
& \quad=H_{n-1} A H_{n-1}-H_{n-1} A H_{n-1}^{\prime}\left(H_{n-1} A H_{n-1}^{\prime}+R_{n-1}\right)^{-1}\left(H_{n-1} A H_{n-1}^{\prime}\right)
\end{align*}
$$

so that

$$
\begin{align*}
& H_{n-1}^{\prime}\left(H_{n-1} A H_{n-1}^{\prime}: R_{n-1}\right) H_{n-1}  \tag{1.31}\\
& \quad=A-A H_{n-1}^{\prime}\left(H_{n-1} A H_{n-1}^{\prime}+R_{n-1}\right)^{-1} H_{n-1} A \\
& \quad=\left(A^{-1}+O_{n}^{-1}\right)^{-1}=\left(A: O_{n}\right)
\end{align*}
$$

by the Schur lemma (see [3], Theorem (8.6)). Hence, (1.1) demonstrates (1.29) in view of (1.30). Again, it is of interest to study $\tau_{1} \cdots \tau_{n}(A)$, however, now for $A \geqq C_{0}>0$.

In order to simplify this study, we introduce the following transformation

$$
\begin{align*}
\sigma_{i}^{*}(A) & =+A+S_{i}  \tag{1.32}\\
\sigma_{i}(A) & =-\phi_{i-1} O_{i-1}\left(A+S_{i}+O_{i-1}\right)^{-1} O_{i-1} \phi_{i-1}^{\prime}
\end{align*}
$$

with $S_{i}=C_{i}+\phi_{i} O_{i} \phi_{i}^{\prime}$.
It will be convenient to rewrite (1.29) as

$$
\begin{equation*}
\tau_{n}(A)=S_{n}-\phi_{n} O_{n}\left(O_{n}+A\right)^{-1} O_{n} \phi_{n} \tag{1.33}
\end{equation*}
$$

Theorem 1.3. For $A \in C$,

$$
\begin{equation*}
\sigma_{i}^{*} \sigma_{i+1} \cdots \sigma_{n}\left(-\phi_{n} O_{n}\left(O_{n}+A\right)^{-1} O_{n} \phi_{n}^{\prime}\right)=\tau_{i} \cdots \tau_{n}(A) \tag{1.34}
\end{equation*}
$$

Proof. We prove this by induction. For $n=i$,

$$
\begin{equation*}
\sigma_{i}^{*}\left(-\phi_{i} O_{i}\left(O_{i}+A\right)^{-1} O_{i} \phi_{i}^{\prime}\right)=S_{i}-\phi_{i} O_{i}\left(O_{i}+A\right)^{-1} O_{i} \phi_{i}^{\prime}=\tau_{i}(A) \tag{1.35}
\end{equation*}
$$

in view of (1.33). Suppose (1.34) is valid for all $n \leqq k$, then

$$
\begin{equation*}
\tau_{i} \cdots \tau_{k+1}(A)=\tau_{i} \cdots \tau_{k}\left(\tau_{k+1}(A)\right) \tag{1.36}
\end{equation*}
$$

which equals by the induction hypothesis

$$
\begin{align*}
\sigma_{i}^{*} & \cdots \sigma_{k}\left(-\phi_{k} O_{k}\left(O_{k}+\tau_{k+1}(A)\right)^{-1} O_{k} \phi_{k}^{\prime}\right)  \tag{1.37}\\
& =\sigma_{i}^{*} \sigma_{i+1} \cdots \sigma_{k+1}\left(\tau_{k+1}(A)-S_{k+1}\right) \\
& =\sigma_{i}^{*} \sigma_{i+1} \cdots \sigma_{k+1}\left(-\phi_{k+1} O_{k+1}\left(O_{k+1}+A\right)^{-1} O_{k+1} \phi_{k+1}^{\prime}\right),
\end{align*}
$$

demgnstrating the assertion for invertible $A$. Since (1.33) holds for all $A \in C$, the theorem is valid.

Corollary 1.5. The maps $\sigma$ and $\tau$ satisfy the relations

$$
\begin{align*}
\sigma_{1}^{*} \sigma_{2} & \cdots \sigma_{n}\left(-\phi_{n} O_{n}\left(O_{n}+C_{n+1}\right)^{-1} O_{n} \phi_{n}^{\prime}\right)  \tag{1.38}\\
& =\tau_{1} \cdots \tau_{n}\left(C_{n+1}\right)=\tau_{1} \cdots \tau_{n+1}(0) \\
& \leqq \tau_{1} \cdots \tau_{n+1}(A) \leqq \tau_{1} \cdots \tau_{n+1}(\infty)=\tau_{1} \cdots \tau_{n}\left(S_{n+1}\right) \\
& =\sigma^{*} \sigma_{2} \cdots \sigma_{n}\left(-\phi_{n} O_{n}\left(O_{n}+S_{n+1}\right)^{-1} O_{n} \phi_{n}^{\prime}\right) .
\end{align*}
$$

Proof. Since $\tau_{1} \cdots \tau_{n+1}(0) \leqq \tau_{1} \cdots \tau_{n+1}(A) \leqq \tau_{1} \cdots \tau_{n+1}(\infty)$, the result follows from the theorem and the obvious relations

$$
\begin{align*}
\tau_{1} \cdots \tau_{n+1}(0) & =\tau_{1} \cdots \tau_{n}\left(C_{n+1}\right),  \tag{1.39}\\
\tau_{1} \cdots \tau_{n+1}(\infty) & =\tau_{1} \cdots \tau_{n}\left(S_{n+1}\right)
\end{align*}
$$

Now for $n \geqq 1$, it is clear that

$$
\begin{equation*}
\sigma_{1}^{*} \sigma_{2} \cdots \sigma_{n+1}(\infty)=\sigma_{1}^{*} \sigma_{2} \cdots \sigma_{n}(0) \equiv A_{n} B_{n}^{-1} \tag{1.40}
\end{equation*}
$$

so that we will suppose that

$$
\begin{equation*}
\sigma_{1}^{*} \sigma_{2} \cdots \sigma_{n+1}(\Gamma)=\left(A_{n+1}+A_{n} K_{n} \Gamma\right)\left(B_{n+1}+B_{n} K_{n} \Gamma\right)^{-1} \tag{1.41}
\end{equation*}
$$

In fact, if the following relations are satisfied

$$
\begin{align*}
K_{n} & =\phi_{n}^{\prime-1} O_{n}^{-1} \\
A_{n+1} & =A_{n} K_{n}\left(S_{n+1}+O_{n}\right)-A_{n-1} K_{n-1} \phi_{n} O_{n}  \tag{1.42}\\
B_{n+1} & =B_{n} K_{n}\left(S_{n+1}+O_{n}\right)-B_{n-1} K_{n-1} \phi_{n} O_{n}
\end{align*}
$$

with $A_{0}=K_{0}^{-1}, A_{1}=S_{1}, B_{0}=O, B_{1}=I$, the relationship (1.41) can be proved by induction. The only point here to be verified is that $B_{n+1}+B_{n} K_{n} \Gamma$ is invertible, and this follows by a simple disconjugacy argument analogous to that given in [3], Theorem (5.1).

Remark 1.3. Equation (1.41) is the analogy of the symplectic system

$$
\left[\begin{array}{l}
X_{n}  \tag{1.43}\\
Y_{n}
\end{array}\right]=\left[\begin{array}{ll}
\phi_{n}^{\prime-1} & \phi_{n}^{\prime-1} O_{n}^{-1} \\
C_{n} \phi_{n}^{\prime-1} & \phi_{n}-C_{n} \phi_{n}^{\prime-1} O_{n}^{-1}
\end{array}\right]\left[\begin{array}{c}
X_{n-1} \\
Y_{n-1}
\end{array}\right]
$$

associated with the mapping $\tau_{k}(\Gamma)$ given by (1.29) as $\tau_{n}(\Gamma)=Y_{n} X_{n}^{-1}$ with $X_{n-1}=I$ and $Y_{n-1}=\Gamma$.

Another point of interest is that of investigation of the equilibrium solutions of (1.33). Previously, we have shown the existence of two equilibrium solutions $L_{i}$ and $V_{i}$. Let $E_{i}$ denote general equilibrium solution, then if $\Sigma_{i}=E_{i}+O_{i-1}$ it follows from (1.33) that

$$
\begin{equation*}
\Sigma_{i}=S_{i}+O_{i-1}-\phi_{i} O_{i}\left(\Sigma_{i+1}\right)^{-1} O_{i} \phi_{i}^{\prime} \tag{1.44}
\end{equation*}
$$

and in the autonomous case

$$
\begin{equation*}
\Sigma=S+O-\phi O(\Sigma)^{-1} O \phi^{\prime} \tag{1.45}
\end{equation*}
$$

These equations generalize the quadratic fixed points of scalar continued fractions.

Example 1.2. If $\Phi=I$ in (1.45), then

$$
\begin{equation*}
\left[O^{-1 / 2} \Sigma O^{-1 / 2}\right]=O^{-1 / 2} C^{-1 / 2}+2 I-\left[O^{-1 / 2} \Sigma O^{-1 / 2}\right]^{-1} \tag{1.46}
\end{equation*}
$$

and

$$
\begin{align*}
O^{-1 / 2} \bar{P} O^{-1 / 2} & =-I+T\left\{1+\frac{1}{2} \mu_{i}+\left[\mu_{i}+\left(\frac{1}{2} \mu_{i}\right)^{2}\right]^{1 / 2}\right\} T^{\prime}  \tag{1.47}\\
& =T\left\{\frac{1}{2} \mu_{i}+\left[\mu_{i}+\left(\frac{1}{2} \mu_{i}\right)^{2}\right]^{1 / 2}\right\} T
\end{align*}
$$

where $T^{\prime} O^{-1 / 2} C O^{-1 / 2} T=\left\{\mu_{i}\right\}$, the notation $\left\{\lambda_{i}\right\}$ denotes the diagonal matrix

$$
\left[\begin{array}{lll}
\lambda_{1} & & 0  \tag{1.48}\\
& \ddots & \\
0 & & \lambda_{n}
\end{array}\right]
$$

The equilibrium solution of $\tau(A)$ is unique and equals $\bar{P}$ with $\bar{P}=L=V$, if $O$ and $C$ are positive definite.

Explicitly, in terms of the unique positive semidefinite square root, we obtain

$$
\begin{equation*}
\bar{P}=\frac{1}{2} C+O^{1 / 2}\left(\frac{1}{4} O^{-1 / 2} C O^{-1} C O^{-1 / 2}+O^{-1 / 2} C O^{-1 / 2}\right)^{1 / 2} O^{1 / 2} \tag{1.49}
\end{equation*}
$$

The continuous version of (1.49) when $C=A \Delta$ and $O=B(1 / \Delta)$ and $\Delta$ tends to zero is

$$
\bar{P}=B^{1 / 2}\left(B^{-1 / 2} A B^{-1 / 2}\right)^{1 / 2} B^{1 / 2}
$$

This equilibrium solution of $\dot{P}=-P B^{-1} P+A$ has been obtained by Reid, Bellman, and others and has lead to the comment that in this particular case the Riccati equation is a continuous "square rooter."

We rewrite (1.29) in the more convenient form for $L_{i}$ and $V_{i}$ as

$$
\begin{equation*}
L_{i}=\phi_{i}^{*} L_{i+1} \phi_{i}^{* \prime}+\phi_{i}\left(L_{i+1}: O_{i}\right) O_{i}^{-1}\left(L_{i+1}: O_{i}\right) \phi_{i}^{\prime}+C_{i} \tag{1.50}
\end{equation*}
$$

with $\phi_{i}^{* \prime}=\left(L_{i+1}+O_{i}\right)^{-1} O_{i} \phi_{i}^{\prime}$ and

$$
\begin{equation*}
V_{i}=\phi_{i}^{* *} V_{i+1} \phi_{i}^{* *}+\phi_{i}\left(V_{i+1}: O_{i}\right) O_{i}^{-1}\left(V_{i+1}: O_{i}\right) \phi_{i}^{\prime}+C_{i} \tag{1.51}
\end{equation*}
$$

with $\phi_{i}^{* * \prime}=\left(V_{i+1}+O_{i}\right)^{-1} O_{i} \phi_{i}^{\prime}$.
It also follows easily that $E_{i}=V_{i}-L_{i}$ satisfies

$$
\begin{equation*}
E_{i}=\phi_{i}^{*} E_{i+1} \phi_{i}^{* * \prime} \tag{1.52}
\end{equation*}
$$

(see [1], Lemma 27). Now the following theorem establishes the asymptotic theory for iterates of $\tau_{i}$.

Theorem 1.4. Suppose $\phi_{i}, R_{i}$ and, $O_{i}$ are invertible and further that there exist $\alpha$ and $\beta$ real positive numbers such that
(i) $0<\alpha I \leqq C_{i}$,
(ii) $V_{i} \leqq \beta I$,
and that $\left\|\phi_{i}-I\right\| \leqq C_{1}$, then $V_{i}=L_{i}=\bar{P}_{i}$ and

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \tau_{i} \cdots \tau_{n}(A)=\bar{P}_{i} \tag{1.53}
\end{equation*}
$$

the convergence being exponential.
Further (1.44) has a unique supnorm bound solution with $\Sigma_{i} \geqq O_{i-1}$ given by $\Sigma_{i}=\bar{P}_{i}+O_{i+1}$ and $\bar{P}_{i} \in C$.

Proof. From (1.50) and (1.51) and our assumptions, it is easily checked that the quadratic forms $\|x\|_{L_{i}}^{2}$ and $\|\mathrm{x}\|_{V_{i}}^{2}$ are Liapounov functions for $\mathbf{x}_{(i+1)}=$ $\phi_{i}^{* \prime} \mathbf{x}_{i}$ and $\mathbf{y}_{(i+1)}=\phi_{i}^{* * *} \mathbf{y}_{i}$, respectively. In particular, this implies

$$
\begin{equation*}
\left\|\phi_{i}^{*} \cdots \phi_{n}^{*}\right\| \leqq c_{*} \exp \left\{-\gamma_{*}(n-i)\right\} \tag{1.54}
\end{equation*}
$$

for $n \geqq i$ for some $c_{*}$ and $\gamma_{*}$ positive and analogous relations for $\phi_{i}^{* * \prime}$ hold with parameters $c_{* *}$ and $\gamma_{* *}$. But since $V_{i}-L_{i}$ satisfies (1.52), it follows that

$$
\begin{equation*}
E_{i}=\phi_{i}^{*} \cdots \phi_{k}^{*} E_{k+1} \phi_{k}^{* *^{\prime}} \cdots \phi_{i}^{* *^{\prime}} \tag{1.55}
\end{equation*}
$$

and hence,

$$
\begin{equation*}
\left\|E_{i}\right\|^{2} \leqq c_{*} c_{* *} \exp \left\{-\left(\gamma_{*}+\gamma_{* *}\right)(k-i)\right\}\left\|E_{k+1}\right\|^{2} \tag{1.56}
\end{equation*}
$$

However, as $k$ is arbitrary and $\left\|E_{k+1}\right\|^{2}$ is uniformly bounded, $E_{i}$ must be zero. The validity of equation (1.53) follows as $L_{i, n} \leqq \tau_{i} \cdots \tau_{n}(A) \leqq V_{i, n}$, since $L_{i, n}$ and $V_{i, n}$ converge to $\bar{P}_{i}=V_{i}=L_{i}$, as $E_{i}=0$. Now suppose (1.44) has two solutions $\Sigma_{i}$ and $T_{i}$, then $x_{i}=\Sigma_{i}-O_{i-1}$ and $y_{i}=T_{i}-O_{i-1}$ are equilibrium solutions of (1.29) both supnorm bounded and members of.C. Now

$$
\begin{align*}
\tau_{i} \cdots \tau_{n}\left(x_{n+1}\right) & =x_{i}  \tag{1.57}\\
\tau_{i} \cdots \tau_{n}\left(y_{n+1}\right) & =y_{i}
\end{align*}
$$

and hence, as $x_{n+1}$ and $y_{n+1}$ are members of $C$,

$$
\begin{align*}
& L_{i, n} \leqq x_{i} \leqq V_{i, n}  \tag{1.58}\\
& L_{i, n} \leqq y_{i} \leqq V_{i, n}
\end{align*}
$$

but the last equation implies $x_{i}=y_{i}=\bar{P}_{i}$.
Remark 1.4. Theorem 1.4 is essentially unchanged when $R_{i}$ is singular and provides stability and a unique equilibrium solution for (1.1) when $L_{i}$ is uniformly positive definite and $V_{i}$ is uniformly bounded in $C$.

Remark 1.5. Notice that conditions (i) and (ii) of the theorem depend in general on the existence of uniform a priori bounds, which we developed previously.

## 2. Applications to the theory of filtering

In order to apply the theory of iterates of the mapping $\tau_{n}(A)$ of the last section to problems in linear filtering theory, we must study the mapping

$$
\begin{equation*}
\Delta_{n}(A)=\phi_{n} A \phi_{n}^{\prime}-K_{n}\left(H_{n-1} A H_{n-1}^{\prime}+R_{n-1}\right)^{\#} K_{n}^{\prime}+G_{n} G_{n}^{\prime} \tag{2.1}
\end{equation*}
$$

with $K_{n}=\phi_{n} A H_{n-1}^{\prime}+G_{n} L_{n-1}^{\prime}, R_{n-1} \geqq L_{n-1} L_{n-1}^{\prime}$ for $A \in C$. This mapping is similar to $\tau_{n}$ except that $H \in M_{s, d}(R)$ and $G \in M_{d, r}(R)$ with $r$ and $s$ for many problems less than $d$. Note also the more general nonlinear term which arises from the observation noise being correlated with the signal process. Iterates of $\Delta$ determine the error covariance matrix of the optimal filter (see [3], Chapters 4 and 9 ). The following lemma relates $\Delta_{n}(A)$ to the simpler Riccati mapping $T_{n}$.

Lemma 2.1. For $A \in C, \Delta_{n}(A)=T_{n}(A)$, where

$$
\begin{equation*}
T_{n}(A)=\psi_{n} S_{n}(A) \psi_{n}^{\prime}+G_{n}\left(I-L_{n-1}^{\prime} R_{n-1}^{\#} L_{n-1}\right) G_{n}^{\prime} \tag{2.2}
\end{equation*}
$$

and

$$
\begin{align*}
S_{n}(A) & =A-A H_{n-1}^{\prime}\left(H_{n-1} A H_{n-1}^{\prime}+R_{n-1}\right)^{\#} H_{n-1} A \\
\psi_{n} & =\phi_{n}-G_{n} L_{n-1}^{\prime} R_{n-1}^{\#} H_{n-1} . \tag{2.3}
\end{align*}
$$

Proof. From the definition, $\Delta_{n}(A)$ can be written as

$$
\begin{align*}
\Delta_{n}(A)=\phi_{n} & S_{n}(A) \phi_{n}^{\prime}-G_{n} L_{n-1}^{\prime}\left(H_{n-1} A H_{n-1}+R_{n-1}\right)^{\#} H_{n-1} A \phi_{n}^{\prime}  \tag{2.4}\\
& \quad+G_{n} L_{n-1}^{\prime}\left(H_{n-1} A H_{n-1}+R_{n-1}\right)^{\#} L_{n-1} G_{n}^{\prime} \\
& \quad-\phi_{n} A H_{n-1}\left(H_{n-1} A H_{n-1}^{\prime}+R_{n-1}\right)^{*} L_{n-1} G_{n}^{\prime}+G_{n} G_{n}^{\prime} .
\end{align*}
$$

But by the definition of $S_{n}(A)$ and Lemma 1.1, it follows that

$$
\begin{align*}
& S_{n}(A) H_{n-1}^{\prime} R_{n-1}^{\#} L_{n-1} G_{n}^{\prime}  \tag{2.5}\\
& \quad=A H_{n-1}\left(H_{n-1} A H_{n-1}^{\prime}+R_{n-1}\right)^{\#} R_{n-1} R_{n-1}^{\#} L_{n-1} G_{n}^{\prime} \\
& \quad=A H_{n-1}\left(H_{n-1} A H_{n-1}^{\prime}+R_{n-1}\right)^{\#} L_{n-1} G_{n}^{\prime},
\end{align*}
$$

and hence,

$$
\begin{align*}
& G_{n} L_{n-1}^{\prime} R_{n-1}^{\#} H_{n-1} S_{n}(A) H_{n-1}^{\prime} R_{n-1}^{\#} L_{n-1} G_{n}^{\prime}  \tag{2.6}\\
& \quad=G_{n} L_{n-1}^{\prime} R_{n-1}^{\#} L_{n-1} G_{n}^{\prime}-G_{n} L_{n-1}^{\prime}\left(H_{n-1} A H_{n-1}+R_{n-1}\right)^{\#} L_{n-1} G_{n}^{\prime} .
\end{align*}
$$

In view of the above equalities and (2.4), the lemma follows.
Remark 2.1. Lemma 2.1 is the discrete time generalization of continuous time equivalence of Riccati equations (see [3], especially page 90 ).

Remark 2.2. Notice that if $R_{n-1}>0$, then $Q$ controllability and $R$ observability of (2.1) hold when and only when (2.2) is $Q^{*}$ controllable and $R^{*}$ observable. In fact, in the general case, it seems appropriate to call (2.1) controllable and observable when these conditions hold for (2.2).

In order to overcome difficulty that $H$ has in general less sensors than the state dimension, we process the observations in blocks of $k$ corresponding to $k$ sequential time observations. In other words, if $\pi\left(n, \Gamma, n_{0}\right)$ represents solution of the Riccati equation, we find the recursion for $\pi\left(n_{0}+k v, \Gamma, n_{0}\right)$ in terms of $\pi\left(n_{0}+k(v-1), \Gamma, n_{0}\right)$. This recursion equation is of the general form of (2.1) with $H_{n} \in M_{k s, d}(R)$, and hence for $s k \geqq d$ using Lemma 2.1, the results of Section 1 are applicable.

As an example, we consider the mapping

$$
\begin{equation*}
T(A)=\phi_{*}\left\{A-A H_{*}^{\prime}\left\{H_{*} A H_{*}^{\prime}+R_{*}\right\}^{*} H_{*} A\right\} \phi_{*}^{\prime}+G_{*} G_{*}^{\prime} . \tag{2.7}
\end{equation*}
$$

Then $T^{k}(A)=\phi^{k}\left\{A-\left(A H^{\prime}+\phi^{-k} G L^{\prime}\right)\left(H A H^{\prime}+R+L L^{\prime}\right)^{\#}\left(H A+L G^{\prime} \phi^{\prime-k}\right) \phi^{\prime k}\right.$ $+G G^{\prime}$, where

$$
L=\left\{\begin{array}{llllll}
0 & 0 & \cdots & & &  \tag{2.8}\\
u_{0} & 0 & \cdots & & & \\
u_{1} & u_{0} & 0 & \cdots & & \\
u_{2} & u_{1} & u_{0} & 0 & \cdots & \\
\cdot & \cdot & \cdot & \cdot & \cdot & \cdot \\
u_{k-2} & u_{k-3} & \cdots & u_{0} & 0 & \cdots
\end{array}\right\}, \quad u_{i}=H_{*} \phi^{i} G_{*},
$$

$$
\begin{gather*}
R=\left[\begin{array}{lllc}
R_{*} & 0 & 0 & 0 \\
0 & R_{*} & 0 & \cdots \\
0 & 0 & R_{*} & 0 \\
\cdot & \cdot & \cdot & \cdot \\
\cdot & \cdot & \cdot & 0 R_{*}
\end{array}\right],  \tag{2.9}\\
H=\left[\begin{array}{c}
H_{*} \\
H_{*} \phi_{*} \\
\vdots \\
H_{*} \phi_{*}^{k-1}
\end{array}\right], \quad G=\left[\phi_{*}^{k-1} G_{*}, \cdots, G_{*}\right] . \tag{2.10}
\end{gather*}
$$

If $\phi_{*}$ is invertible, define $M$ via the equation $H \phi_{*}^{-k} G=M+L$, using properties of the pseudoinverse

$$
\begin{equation*}
\|\mathbf{x}\|_{T^{k}(A)}^{2}=\min _{\mathbf{y} \in R^{s}}\left\{\left\|\psi^{\prime} \mathbf{x}+H^{\prime} \mathbf{y}\right\|_{A}^{2}+\|\mathbf{y}\|_{R+L L^{\prime}}^{2}+\left\|G^{\prime} x\right\|_{I-L^{\prime}\left(R+L L^{\prime}\right)}^{2}{ }_{L}\right\} \tag{2.11}
\end{equation*}
$$

for all $x \in R^{d}$.
An interesting problem is that of characterizing $J_{\ell}=\left\{\underline{x} \in R^{d} \mid T^{j}(A) \underline{x}=T^{\ell}(0) \underline{x}\right.$, all $A \in C$, all $j \geqq \ell\}$. This problem has been solved for $H_{*}$ a vector and $R_{*}=0$ (see [4]). In [5], the general solution has been given for $R=0$. From (2.11) the following theorem determines $J_{k}$ in general.

Theorem 2.1. With matrices given by (2.8), (2.9), and (2.10),

$$
\begin{align*}
J_{k} & =\left\{\mathbf{x} \in R^{d} \mid T^{k}(A) \mathbf{x}=T^{k}(0) \mathbf{x}, \quad \text { for all } A \in C\right\}  \tag{2.12}\\
& =\left\{\mathbf{x} \in R^{d} \mid \psi^{\prime} \mathbf{x}=H^{\prime} \ell, \ell \in R^{s},\|\ell\|_{\left(R+L L^{\prime}\right)}^{2}=0\right\}
\end{align*}
$$

Proof. It is clear that $\left\{\mathbf{x} \in R^{d} \mid T^{k}(A) \mathbf{x}=T^{k}(0) \mathbf{x}\right.$, for all $\left.A \in C\right\} \supseteqq J_{k}$. Since for $A=T^{\ell}(B)$ and $x$ such that $T^{k}(0) \mathbf{x}=T^{k}(A) \mathbf{x}$, it follows that $T^{k}(0) \mathbf{x}=$ $T^{k+\ell}(B) \mathbf{x}$ for arbitrary $B \in C$, so that first set equality is valid. The second set equality follows from (2.11) by considering $T^{k}(0)$ and $T^{k}(A)$ for $A \in \dot{C}$.

Remark 2.3. Notice that the invertibility of $\phi$ is unnecessary for the validity of Theorem 2.1.

The general technique of enlarging the sensor by block processing is valid in the time dependent case and leads to a structure analogous to (2.8), (2.9), and (2.10). Because of this the a priori bounds of Section 1 as well as the asymptotic results apply in general with the only restriction being that there exists a $k$ such that rank $\left[H_{n_{0}}^{\prime}, \cdots \phi^{\prime}\left(n_{0}, n_{0}+k\right) H_{n_{0}+k}^{\prime}\right]=d$ for all $n_{0}$.

## 3. Conclusions

We have shown that the theory of the Riccati equation which arises in the discrete time linear filtering problem can be easily obtained by considering the temporal evolution of $k$ fold iterates. A generalized theory of continued fractions in semidefinite matrices has been given, which provides best possible upper and lower a priori bounds for the Riccati equation solutions. It would
seem that the upper and lower approximates would provide interesting ways to compute suboptimal filters in environments where the prior variance is unknown.

In a future paper, we will study the analogous continuous time situation.

## REFERENCES

[1] W. N. Anderson and R. J. Duffin, "Series and parallel addition of matrices," J. Math. Anal. Appl., Vol. 26 (1969), pp. 576-593.
[2] R. S. Bucy, "Global theory of the Riccati equations," J. Comput. Systems Sci., Vol. 1 (1967), pp. 349-361.
[3] R. S. Bucy and P. D. Joseph, Filtering for Stochastic Processes with Application to Guidance, New York, Interscience, 1968.
[4] R. S. Bucy, D. Rappaport, and L. M. Silverman, "Correlated noise filtering and invariant directions for the Riccati equation," IEEE Trans. Automatic Control, Vol. AC-15 (1970), pp. 535-540.
[5] R. S. Bucy, "Linear and non-linear filtering," Proc. IEEE, Vol. 58 (1970), pp. 854-864.
[6] R. E. Kalman, "A new approach to linear filtering and prediction problems," ASME J. Basic Eng., Vol. 82 (1960), pp. 35-45.
[7] R. E. Kalman and R. S. Bucy. "New results in linear filtering and prediction theory," ASME J. Basic Eng., Vol. 83 (1961), pp. 95-108.
[8] R. E. Kalman and T. S. Englar, "A user's manual for A.S.P.-C," NASA CR-475, Ames Research Center, Moffet Field, California, 1965.
[9] O. Perron, Die Lehre von den Kettenbruecken, Stuttgart, Teubner, 1954.
[10] H. S. Wall, Analytic Theory of Continued Fractions, New York, Chelsea, 1967.

