# A CRITICAL TOPOLOGY IN HARMONIC ANALYSIS ON SEMIGROUPS 

## BY

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## Introduction

Throughout this paper $S$ shall denote a discrete Abelian semi-group with an irreducible unit, denoted 0 , and with a law of cancellation. Spelled out explicitly the two last conditions read:

$$
\begin{align*}
& x_{1}+x_{2}=0 \Rightarrow x_{1}=x_{2}=0,  \tag{1}\\
& x_{1}+y=x_{2}+y \Rightarrow x_{1}=x_{2} \tag{2}
\end{align*}
$$

for elements $x_{1}, x_{2}, y$ belonging to $S$. A semigroup of this kind possesses a natural partial ordering where $x_{1} \leqslant x_{2}$ means that $y \in S$ exists such that $x_{1}+y=x_{2}$. Since $y$ is unique the notation $x_{2}-x_{1}$ stands for an element in $S$ well defined whenever $x_{1} \leqslant x_{2}$.

On $S$ we postulate the existence of a positive function $\omega(x)$, satisfying the following two conditions:

$$
\begin{gather*}
y \leqslant 2 x \Rightarrow \omega(y) \leqslant 2 \omega(x),  \tag{3}\\
\sum e^{-2_{0} \omega(x)} \leqslant 1, \tag{4}
\end{gather*}
$$

where $\lambda_{0}$ is a positive constant. In (4) as in all series in the sequel the summation is extended ower $x \in S$ if no other indication is given. The two previous conditions imply that

$$
\begin{equation*}
N(x, S) \equiv \sum_{y \leqslant x} 1 \leqslant e^{2 \lambda_{0} \omega(x)} \tag{5}
\end{equation*}
$$

The counting function $N(x, S)$ being finite expresses an intrinsic property of $S$ not shared by all semigroups and particularly not by "half planes" of lattice points considered by Helson and Lowdenslager (cf. [3], [4]).

To $S$ and $\omega$ we associate a locally convex topological space of functions $A=A(S, \omega)$ defined as follows: $A$ contains all numeric functions $f(x)$ on $S$ such that for some $\lambda>0$, varying with $f$,

$$
\|f\|_{\lambda}=\sup _{x \in S}|f(x)| e^{-\lambda \omega_{1}(x)}<\infty
$$

Bounded subsets of $A$ are of the form

$$
B_{\lambda, m}=\left\{f \mid\|f\|_{\lambda} \leqslant m\right\}
$$

where $\lambda$ and $m$ are positive constants. A sequence $\left\{f_{n}\right\}_{1}^{\infty}$ converges in $A$ if is contained in a bounded set and converges pointwise, or equivalently expressed, if for some fixed $\lambda,\left\{f_{n}\right\}$ is a Cauchy sequence in the $\lambda$-norm.

Each continuous linear functional on $A$ has the form

$$
(f, \varphi)=\sum f(x) \varphi(x)
$$

where $\varphi$ is a function on $S$ such that for all $\lambda>0$,

$$
\begin{equation*}
\|\varphi\|_{\lambda}^{\prime}=\Sigma|\varphi(x)| e^{\lambda_{\omega}(x)}<\infty . \tag{6}
\end{equation*}
$$

The topology of the dual space $A^{\prime}$ is determined by the family of norms defined by (6).
By $\chi_{A^{\prime}}$ we shall denote all characters on $S$ belonging to $A^{\prime}$, i.e. functions $\xi(x)$ satisfying (6) and the equations

$$
\xi(0)=1, \quad \xi(x+y)=\xi(x) \cdot \xi(y), \quad x, y \in S
$$

The character which equals 1 at $x=0$ and vanishes elsewhere on $S$ does always belong to $\chi_{A^{\prime}}$ and shall be denoted $e$. The Laplace transform of an $f \in A$ is defined by the relation

$$
f(\xi)=(f, \xi)=\sum f(x) \xi(x), \quad \xi \in \chi_{A^{\prime}}
$$

The shift operators $T_{\tau}, \tau \in S$, and their adjoints $T_{-\tau}$ are defined as follows:

$$
\begin{aligned}
& T_{\tau} f(x)=\left\{\begin{array}{cc}
f(x-\tau), & \tau \leqslant x \\
0, & \tau \nLeftarrow x
\end{array}\right. \\
& T_{-\tau} f(x)=f(x+\tau), \quad x, \tau \in S .
\end{aligned}
$$

If $\|f\|_{\lambda}$ is finite it follows by (3) that $\left\|T_{\tau} f\right\|_{2 \lambda} \leqslant\|f\|_{\lambda}$ for all $\tau \in S$, so the set of shift operators is uniformly bounded in $A$. We should also notice that the transform of $T_{\tau} f$ equals $\xi(\tau) f(\xi)$.

After these preliminaries our main problem can be stated. Let $A_{f}$ denote the closed linear subset of $A$ spanned by the set $\left\{T_{\tau} f \mid \tau \in S\right\}$. If $f$ vanishes at a point
$\xi \in \chi_{A^{\prime}}$, then the same is true of the transform of each $g \in A_{f}$. Hence, $A_{f} \neq A$ since $e$ cannot belong to $A_{f}$, the transform $\hat{e}$ being everywhere equal to 1 . The condition $\hat{f} \neq 0$ is thus necessary for $A_{f}$ being the whole of $A$. If this condition is also sufficient we shall say that the closure theorem holds in $A$.

The first and most important result in this general field is Wiener's theorem on the translates in the space $L_{1}(R)$ which subsequently gave rise to the theory of Banach algebras. Without being precise we recall that the closure theorem is known to be true in a variety of Banach spaces provided the topology forces $f(x)$ to tend sufficiently fast to 0 at infinity. This study originates in the belief that the closure theorem would be true again in certain topological spaces on semigroups if the topology admits $f(x)$ to increase sufficiently fast. This conjecture is supported by results in some specific cases. In a previous paper [2] the closure theorem was shown to be false in the Hilbert space $A=L_{2}(S)$ on $S=Z^{+}$(the additive semigroup of integers $\geqslant 0$ ) with norm

$$
\|f\|=\left\{\Sigma|f(x)|^{2}\right\}^{\frac{1}{2}}
$$

However, the necessary and sufficient condition given in [2] and implying the property $A_{f}=A$, permit us to derive the following conclusion: If $\omega(n)=n^{\alpha}, 0<\alpha<\mathrm{I}$, then the closure theorem in $A\left(Z^{+}, \omega\right)$ is false if $\alpha \leqslant 1 / 3$ and true if $\alpha>1 / 2$.

In this paper we aim to show that the validity of the closure theorem in spaces $A(S, \omega)$ depends on a certain critical rate of growth of $\omega$. To characterize that rate of growth shall be our main goal.

## A necessary condition

Theorem I. If the closure theorem holds in $A(S, \omega)$ then the series

$$
\begin{equation*}
\sum_{n=1}^{\infty} \frac{\omega(n x)}{n^{3 / 3}} \tag{7}
\end{equation*}
$$

diverges for each $x \in S, x \neq 0$.
Under the assumption that (7) converges for an $x_{0} \neq 0$ we shall form a counterexample showing that $\hat{f} \neq 0$ does not imply $A_{j}=A$. Let $f_{0}(z)$ be the analytic function

$$
f_{0}(x)=\exp \left(-\frac{1+z}{1-z}\right)=\sum_{0}^{\infty} a_{n} z^{n}, \quad|z|<1
$$

and define $f(x)$ on $S$ by the conditions

$$
f(x)= \begin{cases}a_{n}, & x=n x_{0} \\ 0, & x \notin S_{x_{0}}\end{cases}
$$

where $S_{x_{0}}$ denotes the subsemigroup $\left\{n x_{0}, n=0,1, \ldots\right\}$. This $f$ does certainly belong to $A$ because $f_{0}(z)$ is bounded in the unit disk, and consequently $a_{n}$ is bounded. For each $\xi \in \chi_{A^{\prime}}$, we have $\xi\left(x_{0}\right)=z,|z|<1$. Hence $\xi\left(n x_{0}\right)=z^{n}$ and

$$
f(\xi)=\sum_{0}^{\infty} a_{n} z^{n}=f_{0}(z) \neq 0
$$

In order to prove that $A_{f} \neq A$ it is now sufficient to exhibit an element $\varphi \in A^{\prime}$, $\varphi \equiv 0$, such that for all $\tau \in S$

$$
\begin{equation*}
0=\left(T_{z} f, \varphi\right)=\sum f(x) \varphi(x+\tau) \tag{8}
\end{equation*}
$$

We choose $\varphi$ vanishing outside $S_{x_{0}}$ and equal to $c_{n}$ at $x=n x_{0}$, where the $c_{n}$ shall be determined later. The relation (8) is automatically satisfied for $\tau \notin S_{x_{0}}$ because each term in the series will vanish. If $\varphi \in A^{\prime}$, then the series

$$
\begin{equation*}
\varphi_{0}(z)=\sum_{0}^{\infty} \frac{\varphi(v)}{z^{v}} \tag{9}
\end{equation*}
$$

converges absolutely for $|z| \geqslant 1$ and condition (6) takes the following form for

$$
\begin{gather*}
\tau=n x_{0}, n \geqslant 0 \\
0=\sum_{\nu=0}^{\infty} a_{v} c_{v+n}=\frac{1}{2 \pi} \int_{0}^{2 \pi} f_{0}\left(e^{i \theta}\right) \varphi_{0}\left(e^{i \theta}\right) e^{i n \theta} d \theta \tag{10}
\end{gather*}
$$

where $h_{0}\left(e^{i \theta}\right)=f_{0}\left(e^{i \theta}\right) \varphi_{0}\left(e^{i \theta}\right)$ is a bounded function continuous for $\theta \neq 0(\bmod 2 \pi)$. If therefore (8) is satisfied, then the Fourier coefficients of $h_{0}\left(e^{i \theta}\right)$ vanishes for negative indices and $h_{0}\left(e^{i \theta}\right)$ represents the boundary values of a function $h_{0}(z)$ analytic in $|z|<1$, bounded there, continuous for $|z|=1, z \neq 1$, and vanishing at $z=0$. This implies that $\varphi_{0}(z)$ can be continued analytically into the function $h_{0}(z) / f_{0}(z)$ across each point $\neq 1$ on $|z|=1$. Hence, $\varphi_{0}(z)$ is regular in the region $z \neq 1$. We now recall this classical theorem by Wigert: The series (9) represents a function $\varphi_{0}(z)$ regular for $z \neq 1$ and vanishing at $z=0$, if and only if there exists an entire function $\Phi(w)$ with the properties:

$$
\begin{gather*}
\Phi(v)=c_{v}, \quad v=0,1, \ldots \\
\log |\Phi(w)|=o(|w|), \quad w \rightarrow \infty \tag{11}
\end{gather*}
$$

Wigert's theorem gives us good guidance concerning the choice of the $c_{v}$, but our particular problem needs the following additional result. If (11) is strengthened to

$$
\begin{equation*}
\log |\Phi(w)| \leqslant|w|^{\frac{z}{2}}+\text { const } \tag{12}
\end{equation*}
$$

and if

$$
\begin{equation*}
\sum_{0}^{\infty}|\Phi(v)|<\infty \tag{13}
\end{equation*}
$$

then $\varphi_{0}(z)$ will satisfy the inequality

$$
\begin{equation*}
\left|\varphi_{0}(z)\right| \leqslant \text { const }\left|\exp \left(\frac{1}{8} \frac{1+z}{1-z}\right)\right|, \quad|z|<1 \tag{14}
\end{equation*}
$$

The proof is elementary. The transformation

$$
\psi(s)=\int_{0}^{\infty} e^{-w} \Phi(s w) d w
$$

takes $\Phi$ into an entire function $\psi$ satisfying $\log |\psi(s)| \leqslant|s| / 4+$ const. The inversion formula

$$
\Phi(w)=\frac{1}{2 \pi i} \int_{|s|=r} \psi(s) e^{w / s} \frac{d s}{s}
$$

holds for all finite $w$ and for $r>0$. If the integral representation of $c_{\nu}=\Phi(v)$ is introduced in (9) and the order of integration and summation is reversed we obtain the formula

$$
\varphi_{0}(z)=\frac{z}{2 \pi i} \int_{|s|=r} \frac{\psi(s) d s}{\left(z-e^{1 / s}\right) s},
$$

valid for $z$ outside the contour described by $e^{1 / s},|s|=r$. The estimate on $\psi(s)$ together with an appropriate choice of $r$ yield the inequality

$$
\begin{equation*}
\left|\varphi_{0}(z)\right| \leqslant \text { const } \frac{\exp \left(\frac{1}{4} \frac{1}{|z-1|}\right)}{|z-1|} \tag{15}
\end{equation*}
$$

By virtue of (13), $\varphi_{0}(z)$ is uniformly bounded for $|z|=1, z \neq 1$, and (14) now follows from (15) by an application of the Phragmen-Lindelöf principle to the function

$$
\varphi_{0}(z) \exp \left(-\frac{1}{8} \frac{1+z}{1-z}\right) \quad \text { in }|z|<1
$$

In order to determine $\Phi(w)$ we set $k_{n}=\max \left(\omega\left(v x_{0}\right)\right.$ for $v \leqslant 2 n$. By (3) we have $k_{n} \leqslant 2 \omega\left(n x_{0}\right)$ so the series $\sum n^{-3 / 2} k_{n}$ will converge. Since $k_{n}$ is increasing with $n$ there exists on $(0, \infty)$ a monotonic increasing function $\gamma(u)$ such that $\gamma(n) / k_{n}$ tends to $\infty$ with $n$ while

$$
\begin{equation*}
\int_{1}^{\infty} \frac{\gamma(u)}{u^{3 / 2}} d u=2 \int_{1}^{\infty} \frac{\gamma\left(u^{2}\right)}{u^{2}} d u<\infty \tag{16}
\end{equation*}
$$

We now recall another well-known property of entire functions. If $\gamma(u)$ is in-
creasing and (16) satisfied, then there exists even entire functions $F(w)$ with the properties: $F(0)=1$,

$$
\begin{gather*}
\log |F(w)| \leqslant|w|+\text { const }  \tag{17}\\
|F(u)| \leqslant \text { const } e^{-\gamma\left(u^{\prime}\right)}, \quad u \text { real. } \tag{18}
\end{gather*}
$$

We choose such a function $F(w)$ and define $\Phi(w)=F(\sqrt{w})$. Then $\Phi(w)$ is entire and satisfies (12). Moreover,

$$
|\Phi(v)| \leqslant \text { const } e^{-\gamma(\nu)} \leqslant \text { const } e^{-\lambda_{v} \omega\left(\nu x_{0}\right)}
$$

where $\lambda_{\nu}$ tends to $\infty$ as $\nu \rightarrow \infty$. The choice $c_{\nu}=\Phi(\nu)$ therefore yields a function $\varphi \in A^{\prime}$ and the condition (10) is certainly satisfied since the product $\varphi_{0}(z) f_{0}(z)$ is a bounded analytic function in $|z|<1$, vanishing at $z=0$. Hence, $A_{f} \neq A$.

## Multipliers

Assume $f, g \in A,\|f\|_{\lambda}<\infty,\|g\|_{\lambda^{\prime}}<\infty$, then the pointwise product $f \cdot g$ belongs to $A$ and $\|f \cdot g\|_{\lambda+\lambda^{\prime}} \leqslant\|f\|_{\lambda}\|g\|_{\lambda^{\prime}}$. The convolution $f * g$ is defined by the formula

$$
f * g(x)=\sum_{x_{1}+x_{2}-x} f\left(x_{1}\right) g\left(x_{2}\right)=\sum_{y \leqslant x} f(y) g(x-y) .
$$

By virtue of (3), $\quad|f(y) g(x-y)| \leqslant\|f\|_{\lambda}\|g\|_{\lambda \cdot} e^{2\left(\lambda+\lambda^{\prime}\right) \omega(x)}$,
which together with (5) implies

$$
\|f * g\|_{x^{\prime \prime}} \leqslant\|f\|_{\lambda}\|g\|_{x^{2}}
$$

for $\lambda^{\prime \prime} \geqslant 2\left(\lambda+\lambda^{\prime}+\lambda_{0}\right)$. Both pointwise multiplication and convolution are therefore continuous mappings of $A \times A$ into $A$.

Of particular interest is the properties of $A$ considered as a convolution algebra. By $J(f)$ we shall denote the ideal generated by $f$ :

$$
J(f)=\{f * g \mid g \in A\}
$$

and by $\bar{J}(f)$ the closure of $J(f)$. Since each linear combination of the elements $T_{\tau} f$, $\tau \in S$, equals a convolution $f * k$ where $k$ has finite support, it follows that $A_{f}=\breve{J}(f)$. We now introduce the following notion: a function $\varrho \in A$ shall be called a converting multiplier if for each $f \in A$ with $\hat{f} \neq 0$ it holds that $\varrho \cdot f_{1} \in \bar{J}(f)$ whenever $f_{1} \in \bar{J}(f)$.

The collection $M$ of converting multipliers is obviously a closed linear subset of $A$ with the property that $\varrho_{1}, \varrho_{2} \in M$ implies $\varrho_{1} \cdot \varrho_{2} \in M$. This notion is connected with our main problem as follows: The closure theorem holds in $A$ if and only if $M=A$.

The proof is trivial. If the closure theorem holds, then for each $f$ with $\hat{f} \neq 0$ and for each $\varrho \in A$ we have $\varrho \cdot f_{1} \in A=\bar{J}(f)$ and the conclusion $\varrho \in M$ follows. If, on the other hand $M=A$, then $M$ contains the unit $e$ of the convolution algebra, and $\hat{f} \neq 0$ implies $e \cdot f \in \bar{J}(f)$. But $e \cdot f=e \cdot f(0)$ and $f(0)=\hat{f}(e) \neq 0$, so $\bar{J}(f)$ contains $e$ and is therefore equal to the whole of $A$.

The notion of converting multiplier is thus trivial whenever the closure theorem holds in the space. This is however not the case with a subset $M_{0}$ of $M$ defined as follows: An element $\varrho \in A$ is a proper converting multiplier and shall belong to $M_{0}$ if for each $f \in A$ with $\hat{f} \neq 0$ and for each $g \in A$ the relation $\varrho \cdot(f * g)=f * k$ is satisfied by some element $k \in A$.

The closure $\bar{M}_{0}$ of $M_{0}$ is contained in $M$, and $M_{0}$ contains finite products of its elements. The question whether always $\bar{M}_{0}=M$ has not been resolved in this paper.

The set $M_{0}$ derives its importance from the fact that it contains subsets which can be derived by a simple algebraic method, as will be shown in the following section.

## Polynomials on $S$

By $H=H(S)$ we shall denote the set of all additive (and finite) functions $\theta(x)$ on $S$. Each mapping $x \rightarrow \theta(x)$ is thus an homomorphism of $S$ into an additive semigroup of complex numbers. A function $p(x)$ shall be called a polynomial on $S$ if it has a representation

$$
\begin{equation*}
p(x)=p(0)+\sum_{n} \prod_{m} \theta_{m . n}(x) \tag{19}
\end{equation*}
$$

where the $\theta_{m, n}$ belong to $H$ and where series and products are finite. In (19), $p(0)$ stands for the function equal to the constant $p(0)$ everywhere on $S$.

We shall first derive some properties valid, irrespective of the topology, for all functions $f(x)$ which are finite on $S$. Since the number $N(x, S)$ of elements $y \leqslant x$ is finite, it follows that the convolution $f * g$ is always well defined. We shall write $h^{* n}$ for the $n$-fold convolution of $h$ with itself, defined as $e$ for $n=0$. The value of $h^{* n}$ at a point $x$ does not change if $h$ is replaced by the function $h_{0}(y)$ which equals $h(y)$ for $y \leqslant x$ and vanishes elsewhere on $S$. For $h_{0}$ we have the familiar inequality

$$
\Sigma\left|h_{0}^{* n}(y)\right| \leqslant\left\{\sum\left|h_{0}(y)\right|\right\}^{n} .
$$

Consequently

$$
\left|h^{* n}(x)\right|=\left|h_{0}^{* n}(x)\right| \leqslant\left\{\sum_{y \leqslant x}|h(y)|\right\}^{n}
$$

If therefore $h(x)$ is finite on $S$, the series

$$
\begin{equation*}
f(x)=\sum_{0}^{\infty} \frac{h^{* n}(x)}{n!} \tag{20}
\end{equation*}
$$

will always converge absolutely. We should also notice that if $h(0)=0$, then for fixed $x, h^{* n}(x)$ will vanish for all $n$ sufficiently large. This follows from the fact that the equation $\sum_{1}^{n} x_{\nu}=x, x_{\nu} \neq 0$, has no solution if $n>(N(x, S)-1)^{2}$.

In order to avoid any confusion we denote by $F$ the set of all finite numeric functions on $S$, by $E$ the set of functions representable by the series (20) with $h \in F$, and by $P$ the set of polynomials on $S$. The following lemma will play an important role in this study:

Lemma I. Assume $p \in P, f \in E$ and $g \in F$. Then the relation

$$
\begin{equation*}
p \cdot(f * g)=f * k \tag{21}
\end{equation*}
$$

is always satisfied by some element $k \in F$.
If $\theta \in H$, then the value of $\theta \cdot(g * h)$ at a point $x$ can be written

$$
\sum_{y \leqslant x}\{\theta(y) g(y) h(x-y)+g(y) \theta(x-y) h(x-y)\} .
$$

Consequently

$$
\begin{equation*}
\theta \cdot(g * h)=(\theta \cdot g) * h+g *(\theta \cdot h) \tag{22}
\end{equation*}
$$

By iteration of this formula we obtain

$$
\begin{equation*}
\theta \cdot h^{* n}=n h^{* n-1} *(\theta \cdot h) \tag{23}
\end{equation*}
$$

It therefore (20) is multiplied by $\theta$ if follows that

$$
\begin{equation*}
\theta \cdot f=\sum_{n=0}^{\infty} \frac{n h^{* n-1}}{n!} *(\theta \cdot h)=f *(\theta \cdot h) \tag{24}
\end{equation*}
$$

Another application of (22) yields the more general formula

$$
\begin{equation*}
\theta \cdot(f * g)=f *\{\theta \cdot g+g *(\theta \cdot h)\} . \tag{25}
\end{equation*}
$$

For $h$ and $f$ fixed we denote by $U_{\theta}$ the linear operator: $g \rightarrow \theta \cdot g+g *(\theta \cdot h)$. If $\left\{\theta_{\nu}\right\}_{1}^{q}$ is a finite sequence $\in H$, then the relation

$$
\begin{equation*}
\prod_{1}^{0} \theta_{\nu} \cdot(f * g)=f * k \tag{26}
\end{equation*}
$$

is satisfied by $k=U_{\theta_{1}} U_{\theta_{3}} \ldots U_{\theta_{q}} g \in F$. Therefore (21) has a solution $k=p(0) g+\sum k_{n}$, where the $k_{n}$ satisfy equations of the form (26).

We now return to the space $A$ and denote by $P_{A}$ the subset of polynomials generated by functions $\theta \in H \cap A$. An $f \in A$ possessing a representation (20) with $h \in A$ shall be called exponential. As a consequence of this definition we shall have

$$
\hat{f}(\xi)=\sum_{0}^{\infty} \frac{\hat{h}(\xi)}{n!}=e^{\hat{h}(\xi)}, \quad \xi \in \chi_{A^{\prime}},
$$

so $\hat{f} \neq 0$ is a prerequisite for $f$ being exponential.
Since $A$ is an algebra both under multiplication and convolution, the operators denoted $U_{\theta}$ are bounded in $A$ whenever $h$ and $\theta$ belong to $A$, and Lemma I thus asserts that $p \cdot(f * g)$ belongs to $J(f)$ if $f$ is exponential and $p \in P_{A}$.

We can now summarize: If $\hat{f} \neq 0$ implies that $f$ is exponential then all polynomials $\epsilon P_{A}$ are proper converting multipliers and the closure theorem holds in $A$ if $e$ is contained in the closure of $P_{A}$. The original problem has herewith branched out into two separate questions.

## Exponential elements in $A$

Conclusive results on our main problem requires further information about $S$ and $\omega$. This should be obvious already by the fact that the conditions introduced so far do not imply that $\chi_{A}$. contains any other character than $\xi=e$. In order to remedy this situation we observe that the topology of $A$ remains unchanged if $\omega$ is replaced by a function $\omega_{1}$ which is equivalent with $\omega$ in the sense that

$$
\begin{equation*}
k^{-1} \leqslant \frac{\omega_{1}(x)}{\omega(x)} \leqslant k \tag{27}
\end{equation*}
$$

for some constant $k>1$. Of particular significance for our problem is the subset $H^{+}$ of $H$ consisting of real valued additive functions $\vartheta(x)$ tending to $+\infty$ as $x \rightarrow \infty$ in $S$. Such a function is obviously strictly positive for $x \neq 0$. Our new condition reads: $H^{+}$ contains an element $\vartheta(x)$ such that $\omega(x)$ is equivalent with a function of the form $\psi(\vartheta(x))$, where $\psi(r)$ is positive and increasing for $r \geqslant 0$ with growth limited by the inequalities

$$
\begin{equation*}
c_{1} \log r \leqslant \psi(r)=o(r), \quad r \rightarrow \infty . \tag{28}
\end{equation*}
$$

The first inequality implies $\vartheta \in A$, and the second together with (4) imply that the character $e^{-s \vartheta(x)}$ belongs to $\chi_{A^{\prime}}$ if $s$ is a complex number with positive real part.

Lemma II. Let $S$ and $\omega$ satisty the previously stated conditions. Then each $f \in A$ with non-vanishing transform is exponential.

Let us first show that there exist constants $k_{1}$ and $k_{2}$ such that for $x, y \in S$.

$$
\begin{equation*}
\frac{\omega(y)}{\omega(x)} \leqslant k_{1} \frac{\vartheta(y)}{\vartheta(x)}+k_{2} . \tag{29}
\end{equation*}
$$

By virtue of (27) the inequality is satisfied for $\vartheta(y) \leqslant \vartheta(x)$ if $k_{2} \geqslant k^{2}$. If $\vartheta(y)>\vartheta(x)$ we set $\vartheta(y)=r, \vartheta(x)=r_{0}$ and define $n \geqslant 1$ so that $2^{n-1} r_{0}<r \leqslant 2^{n} r_{0}$. Consequently

$$
\frac{\omega(y)}{\omega(x)} \leqslant k^{2} \frac{\psi\left(2^{n} r_{0}\right)}{\psi\left(r_{0}\right)} .
$$

By (3) we have $\omega\left(2^{n} x\right) \leqslant 2^{n} \omega(x)$. Hence, $\psi\left(2^{n} r_{0}\right) / \psi\left(r_{0}\right) \leqslant k^{2} 2^{n}$, and (29) is satisfied if we choose $k_{1} \geqslant 2 k^{4}$.

Let $G$ be the minimal extension of $S$ to a group and let $\hat{G}$ be the compact Abelian group which is the dual of $G$. By $\beta=\beta(x)$ we denote characters on $G$ of modulus 1 , and by $d \beta$ Haar's measure on $\hat{G}$ normalized by the condition $\int d \beta=1$. We shall use the notation

$$
f_{s}(x)=e^{-s \vartheta(x)} f(x) .
$$

Each mapping $f \rightarrow f_{s}$ takes $A$ into the space $L_{1}(S)$ with norm $\|f\|_{L_{1}}=\sum|f(x)|$. For the $n$-fold convolution of $f_{s}$ we have

$$
\begin{equation*}
f_{s}^{* n}(x)=e^{-s \theta(x)} f^{* n}(x) \tag{30}
\end{equation*}
$$

Without loss of generality we may assume $f(0)=1$ and write $f=e+g$, with $g(0)=0$. Let $\sigma_{0}$ be so large that $\left\|g_{s}\right\|_{L_{1}} \leqslant \frac{1}{2}$ for $s=\sigma+i t, \sigma \geqslant \sigma_{0}$. Then

$$
f\left(e^{-s \vartheta} \beta\right)=1+\sum g_{s}(x) \beta(x)=1+\hat{g}_{s}(\beta)
$$

where $\left|g_{s}(\beta)\right| \leqslant \frac{1}{2}$ for each $\beta$. The logarithm of this function is now uniquely determined on $\hat{G}$ by the formula

$$
\begin{equation*}
\log f\left(e^{-s \vartheta} \beta\right)=\sum_{1}^{\infty} \frac{(-1)^{n+1}}{n} \hat{g}_{s}^{n}(\beta), \quad \sigma \geqslant \sigma_{0} \tag{31}
\end{equation*}
$$

where

$$
\hat{g}_{s}^{n}(\beta)=\sum g_{s}^{* n}(x) \beta(x)
$$

Since $g(0)=0$ we know that $g^{* n}(x)=0$ for $x$ fixed if $n$ is sufficiently large. A function $h(x)$ is therefore well defined on $S$ by the relation

$$
\begin{equation*}
\sum_{1}^{\infty} \frac{(-1)^{n+1}}{n} g_{s}^{* n}(x)=e^{-s \theta(x)} \sum_{1}^{\infty} \frac{(-1)^{n+1}}{n} g^{* n}(x)=e^{-s \vartheta(x)} h(x), \tag{32}
\end{equation*}
$$

and $h(0)=0$. The left-hand side of (31) is a continuous function on $\hat{G}$, depending on a parameter $s$, with an absolutely convergent Fourier series:

$$
\begin{equation*}
\log \hat{f}\left(e^{-s \vartheta} \beta\right)=\sum_{x \in G} e^{-s \vartheta(x)} h(x) \beta(x), \quad \sigma \geqslant \sigma_{0} \tag{33}
\end{equation*}
$$

where $h$ is defined $=0$ outside $S$. Hence

$$
\begin{equation*}
e^{-s \theta(x)} h(x)=\int_{\hat{G}} \log f\left(e^{-s \vartheta} \beta\right) \bar{\beta}(x) d \beta, \quad \sigma \geqslant \sigma_{0} . \tag{34}
\end{equation*}
$$

Since $\hat{f} \neq 0$ in $\chi_{A}$, the logarithm has a unique analytic extension to the whole right half plane and the formula (34) holds there by analytic continuation. In particular, since $h(0)=0$,

$$
\begin{equation*}
0=\int_{\hat{G}} \log \hat{f}\left(e^{-s \theta} \beta\right) d \beta \tag{35}
\end{equation*}
$$

Hence,

$$
\begin{equation*}
e^{-\sigma \theta(x)}|h(x)| \leqslant 2 \max _{\beta} \log \left|\hat{f}\left(e^{-\sigma \vartheta} \beta\right)\right|, \quad \sigma>0 . \tag{36}
\end{equation*}
$$

Let $\lambda$ be so large that $\|f\|_{\lambda} \leqslant 1$. Then

$$
\left|\hat{f}\left(e^{-\sigma \theta} \beta\right)\right| \leqslant \sum e^{-\sigma \theta(x)+\lambda \omega(x)} .
$$

On defining

$$
m(\varepsilon)=\sup _{y \in S}(-\varepsilon \vartheta(y)+\omega(y))
$$

we obtain by virtue of (4)

$$
\log \left|\hat{f}\left(e^{-\sigma \vartheta} \beta\right)\right| \leqslant\left(\lambda+\lambda_{0}\right) m\left(\frac{\sigma}{\lambda+\lambda_{0}}\right)
$$

Hence, for all $\sigma>0, \quad|h(x)| \leqslant 2\left(\lambda+\lambda_{0}\right) m\left(\frac{\sigma}{\lambda+\lambda_{0}}\right) e^{\alpha \theta(x)}$.
In this relation we choose $\sigma=k_{1}\left(\lambda+\lambda_{0}\right) \omega(x) / \vartheta(x)$, where $k_{1}$ is the constant occurring in (29). We want to show that this choice yields

$$
|h(x)| \leqslant c_{1} e^{\lambda_{1} \omega(x)}
$$

with $c_{1}$ independent of $x$ and with $\lambda_{1}=k_{1}\left(\lambda+\lambda_{0}\right)+1$. We have thus to show that for $x, y \in S$.

$$
2\left(\lambda+\lambda_{0}\right)\left\{-k_{1} \frac{\vartheta(y) \omega(x)}{\vartheta(x)}+\omega(y)\right\} \leqslant c_{1} e^{\omega(x)}
$$

or equivalently that
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$$
\frac{\omega(y)}{\omega(x)} \leqslant k_{1} \frac{\vartheta(y)}{\vartheta(x)}+\frac{c_{1}}{2\left(\lambda+\lambda_{0}\right)} \cdot \frac{e^{\omega(x}}{\omega(x)}
$$

and this inequality is satisfied by virtue of (29) if $c_{1}$ is chosen properly. We have thus shown that $h \in A$.

By the definition of $h$ we have for $\sigma \geqslant \sigma_{0}$,

$$
\begin{gathered}
f\left(e^{-s \vartheta} \beta\right)=\exp \left\{\sum_{x \in S} h_{s}(x) \beta(x)\right\} \\
f_{s}(x)=\sum_{0}^{\infty} \frac{h_{s}^{* n}(x)}{n!}
\end{gathered}
$$

implying

In view of (30) this yields the requested representation

$$
f(x)=\sum_{0}^{\infty} \frac{h^{* n}(x)}{n!}
$$

It is well known from the elementary theory of Taylor series that the sole condition that $\omega$ is monotonic increasing is not sufficient in order to imply $f$ exponential if $\hat{f} \neq 0$. Some additional property is required preventing the rate of growth of $\omega$ to change too erratically. We want to point out that the lemma remains true under the assumption that there exists a $\vartheta \in H^{+}$and a positive constant $\varepsilon$ such that for $x$ outside some finite set,

$$
\vartheta^{\varepsilon}(x)<\omega(x)<\vartheta^{1-\varepsilon}(x) .
$$

## Polynomial approximation of $\boldsymbol{e}$

Theorem II. Let $S$ and $\omega$ satisfy the previously introduced condition, and assume in addition that $\psi(r)$ in (28) is a convex function of $\log r$. Then the following is true: The closure of the set of polynomials $P_{A}$ contains the unit e and the closure theorem holds in $A$ if and only it the divergence condition of Theorem I is satisfied.

If $\vartheta \in H^{+}$then the mapping $x \rightarrow \vartheta(x)$ takes $S$ to a discrete set of numbers $\geqslant 0$, and $\vartheta(x)$ has a positive minimum $r_{0}$ for $x \neq 0$. In order to prove that $e$ is contained in the closure of $P_{A}$ it is thus sufficient to show the existence of a sequence of polynomials $Q_{n}(t)$ assuming the value 1 at $t=0$ and such that $Q_{n}(t) e^{-\varphi(t)}$ converges uniformly to 0 for $t \geqslant r_{0}$. Then $p_{n}(x)$ defined as $Q_{n}(\vartheta(x))$ will belong to $P_{A}$ and converge to $e$ in the space $A$. The existence of $Q_{n}$ can be considered as a special case of Bernstein's classical approximation problem, formulated for the positive real axis
$R^{+}$. Let $C_{0}\left(R^{+}\right)$denote the space of functions continuous for $t \geqslant 0$ and vanishing at $\infty$, and let $\psi(t)$ be a given function continuous on $R^{+}$and tending so fast to $+\infty$ that

$$
\begin{equation*}
A_{n}=\sup _{t \geqslant 0} t^{n} e^{-\psi(t)}<\infty, \quad n=0,1,2 \ldots \tag{37}
\end{equation*}
$$

The problem is to decide whether linear combinations of $t^{n} e^{-\psi(t)}, n=0,1,2 \ldots$, are dense in $C_{0}\left(R^{+}\right)$. Bernstein's original results imply that approximation is possible if $e^{\nu(t)}$ has a minorant for $t \geqslant 0$ of the form

$$
\begin{gather*}
F(t)=\sum_{0}^{\infty} c_{v} t^{\nu}, \quad c_{0}>0, c_{v} \geqslant 0 \\
\int_{1}^{\infty} \frac{\log F(t)}{t^{3 / 2}} d t=\infty \tag{38}
\end{gather*}
$$

This result applies immediately to the problem at hand. The divergence condition in Theorem I implies

$$
\begin{equation*}
\int_{1}^{\infty} \frac{\psi(t)}{t^{3 / 3}} d t=\infty \tag{39}
\end{equation*}
$$

where $\psi$ is the function in condition (28). Moreover, if $\psi(t)$ is a convex function of $\log t$, then (39) implies

$$
\begin{equation*}
\lim _{t=\infty} \frac{\psi(t)}{\log t}=\infty \tag{40}
\end{equation*}
$$

so (37) is satisfied. In order to show existence of minorants $F(t)$ it suffices to choose

$$
F(t)=\sum_{0}^{\infty} \frac{t^{n}}{2^{n} A_{n}}
$$

where $A_{n}$ is defined by (37). Due to the convexity of $\psi$ there exists for each $n \geqslant 0$ a number $t_{n}$ such that $e^{\psi\left(t_{n}\right)}=t_{n}^{n} / A_{n}$. Hence, $\log F\left(2 t_{n}\right)>\psi\left(t_{n}\right)$ and a simple computation shows that (39) implies (38). A sequence $Q_{n}(t)$ with the requested properties does therefore exist since any continuous function equal to 1 at $t=0$ and vanishing for $t \geqslant r_{0}$, can be approached uniformly on $R^{+}$by functions $Q_{n}(t) e^{-\psi(t)}$. This finishes: the proof of Theorem II since we already know that convergence in (7) implies that: the closure theorem is false and consequently $e$ not contained in the closure of $P_{A}$.

It should be pointed out that without the additional convexity condition the preceding analysis does not imply that the closure theorem is false in $A$ if $e \ddagger \bar{P}_{A}$. This problem remains unsolved even in the case $S=Z^{+}$.

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