

On the Carleman determinants

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Fredholm [1] defined the notion of determinants and subdeterminants of all orders for the integral equation

$$(1) \quad x(s) + \int_0^1 \tau(s, t) x(t) dt = x_0(s)$$

with continuous kernel $\tau(s, t)$. Fredholm's determinant theory was recently generalized (Grothendieck [1-2], Leżański [1-2], Ruston [1-2], Sikorski [1-4]; see also the expository paper by Sikorski [5]) over the case of linear equations

$$(2) \quad x + Tx = x_0$$

in arbitrary Banach spaces X , under some hypothesis on the operators T .

In the general theory of determinants in Banach spaces X , two notions of the subdeterminant of an order n can be introduced: the first one coincides with the algebraic notion of subdeterminants in the case where X is finitely dimensional; the second one coincides with the original Fredholm subdeterminant in the case where X is the space \mathcal{O} (for details, see Grothendieck [1] and Sikorski [4-5]).

In the case where X is a Hilbert space, the general determinant theory is applicable to (2) if and only if T is *nuclear*, i. e. (see e. g. Sikorski [2], Theorem V) if $T = T_1 T_2$, where T_1, T_2 are Hilbert-Schmidt operators. An operator T in a Hilbert space X is said to be a *Hilbert-Schmidt operator* if, given any orthogonal system of coordinates in X (i. e. a complete orthonormal set in X), T is represented by an infinite square matrix $(\tau_{i,j})$ such that

$$(3) \quad \|T\| = \sqrt{\sum_{i,j} |\tau_{i,j}|^2} < \infty.$$

In the case of the Hilbert space L^2 , Carleman [1] (see also Hille and Tamarkin [1-2], Smithies [2]) defined the notion of determinant and subdeterminants for the integral equation (1) under the hypothesis that

$$\int_0^1 \int_0^1 |\tau(s, t)|^2 ds dt < \infty.$$

Carleman's formulae were a simple modification of those of Fredholm: everywhere the expressions $\int_0^1 \tau(s, s) ds$ were replaced by zero, following an idea of Hilbert [1].

In the abstract formulation (2), the case investigated by Carleman is that where T is a Hilbert-Schmidt operator. A very interesting treatment of the Carleman determinant and the first subdeterminant in an abstract Hilbert space X is due to Smithies [1] who also gave a simple proof of the convergence of series defining the determinant and the first subdeterminant ⁽¹⁾.

The subject of this paper is to give a complete determinant theory for the equation (2) in an abstract Hilbert space X , T being a Hilbert-Schmidt operator. The subdeterminants of all orders are defined. Similarly, as in the general determinant theory in Banach spaces, two notions of subdeterminants of an order n are introduced. They are called, respectively, the Carleman subdeterminant and the Carleman-Fredholm subdeterminant. Carleman's original subdeterminants coincide with the Carleman-Fredholm subdeterminants, in the terminology assumed in this paper.

It is easy to know what formulae should define the subdeterminants of (2). Indeed, the subdeterminants should be defined by the same series as in the case of a nuclear T , with only one modification: everywhere the trace of T (i. e. the abstract substitute of $\int_0^1 \tau(s, s) ds$) should be replaced by 0. The only difficulty is to prove that the series converge. Fortunately, this difficulty is only on the surface. Using an argument of Grothendieck [1], the convergence of series defining the subdeterminants can be deduced from the Carleman-Smithies theorem stating the convergence of the series defining the determinant.

§ 1. Terminology and notation. We shall consider a fixed Hilbert space X . The letter \mathcal{E} will denote the Hilbert space of all linear bounded functionals on X . The letters x, y, z (with indices, if necessary) denote always elements of X , and the letters ξ, η, ζ — elements of \mathcal{E} . The value of a functional ξ at a point x is denoted by ξx .

The symbol \mathfrak{D}_n will denote the Banach space of all $2n$ -linear bounded functionals on $\mathcal{E}^n \times X^n$. If $B \in \mathfrak{D}_n$, then the value of B at a point $(\xi_1, \dots, \xi_n, x_1, \dots, x_n) \in \mathcal{E}^n \times X^n$ will be usually denoted by

$$B \begin{pmatrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{pmatrix}.$$

⁽¹⁾ For another theory of determinants of operators in an abstract Hilbert space, see Fuglede and Kadison [1, 2].

The norm of B in \mathfrak{D}_n is

$$|B| = \sup_{|\xi_1| \leq 1, \dots, |\xi_n| \leq 1, |x_1| \leq 1, \dots, |x_n| \leq 1} \left| B \begin{pmatrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{pmatrix} \right|.$$

In particular, \mathfrak{D}_1 is the Banach space of all bounded bilinear functionals A on $\mathcal{E} \times X$. If $A \in \mathfrak{D}_1$, we write also ξAx instead of $A \begin{pmatrix} \xi \\ x \end{pmatrix}$. Elements in \mathfrak{D}_1 are called *operators*.

Let $A \in \mathfrak{D}_1$. For every x there exists exactly one y such that $\xi y = \xi Ax$ for every ξ . We denote this element y by Ax . Similarly, for every ξ there exists exactly one η such that $\eta x = \eta Ax$ for every x . We denote this element η by ξA . Obviously, the mappings $y = Ax$ and $\eta = \xi A$ are adjoint endomorphisms in X and \mathcal{E} respectively, and their ordinary norms coincide with the norm $|A|$ of A in \mathfrak{D}_1 . Conversely, every endomorphism in X or in \mathcal{E} is of the above form. Thus every operator can be simultaneously interpreted as a bilinear functional on $\mathcal{E} \times X$, or an endomorphism in X , or as an endomorphism in \mathcal{E} . The three interpretations of any $A \in \mathfrak{D}_1$ will be systematically used in this paper.

The set \mathfrak{D}_1 of operators is a Banach algebra with the following definition of the product $A_1 A_2$ of $A_1, A_2 \in \mathfrak{D}_1$:

$$\xi(A_1 A_2)x = (\xi A_1)(A_2 x).$$

In other words, the product $A_1 A_2$ interpreted as an endomorphism in X (in \mathcal{E}) is the superposition of the endomorphism A_1, A_2 in X (of the endomorphism A_2, A_1 in \mathcal{E}). The unit element of the algebra \mathfrak{D}_1 is the fundamental bilinear functional I :

$$\xi I x = \xi x.$$

By definition, $I x = x$ and $\xi I = \xi$ for all x and ξ .

Let x_0, ξ_0 be fixed. The operator K_0 defined by the formula

$$\xi K_0 x = \xi x_0 \cdot \xi_0 x$$

(i. e. the product of numbers ξx_0 and $\xi_0 x$) is called *one-dimensional* and denoted by $x_0 \cdot \xi_0$. By definition, $K_0 x = x_0 \cdot \xi_0 x$, and $\xi K = \xi x_0 \cdot \xi_0$ (the dot replaces here parentheses).

Any finite sum $K = \sum_{i=1}^m x_i \cdot \xi_i$ of one-dimensional operators is called *finitely dimensional operator*.

The letter T denotes always Hilbert-Schmidt operators, and the letter \mathfrak{S} denotes the Banach algebra of all Hilbert-Schmidt operators T with the norm $\|T\|$ defined by (3). \mathfrak{S} is an ideal in \mathfrak{D}_1 . Moreover,

$$(4) \quad \|x_0 \cdot \xi_0\| = |x_0| \cdot |\xi_0|$$

$$(4') \quad |T| \leq \|T\| \text{ for every } T \in \mathfrak{S},$$

and

$$(4'') \quad \|TA\| \leq \|T\| \cdot \|A\|, \|AT\| \leq \|A\| \cdot \|T\| \text{ for every } A \in \mathfrak{D}_1 \text{ and } T \in \mathfrak{S}.$$

Any nuclear operator T , i. e. a product $T = T_1 T_2$ of two Schmidt operators T_1, T_2 , has a well defined trace denoted by $\text{tr} T$. Viz. if, in a given system of orthogonal coordinates in X , T_1 and T_2 are represented by matrices $(\tau_{i,j})$ and $(\sigma_{i,j})$ respectively, then

$$\text{tr} T = \sum_{i,j} \tau_{j,i} \sigma_{i,j}.$$

The number $\text{tr} T$ just defined depends neither on the representation of T in the form $T_1 T_2$, nor on the choice of the system of coordinates in X . We have

$$(5) \quad \text{tr} ST = \text{tr} TS \quad \text{and} \quad |\text{tr} ST| \leq \|S\| \cdot \|T\|$$

for $T, S \in \mathfrak{S}$.

Suppose that $B(\xi_1, \dots, \xi_n)$ is a $2n$ -linear functional on $\mathcal{E}^n \times X^n$ such that B considered as function of ξ_n and x_n only is a Hilbert-Schmidt operator, i. e., for any fixed $\xi_1, \dots, \xi_{n-1}, x_1, \dots, x_{n-1}$,

$$B(\xi_1, \dots, \xi_n) = \xi_n S x_n,$$

where $S \in \mathfrak{S}$. Let $T \in \mathfrak{S}$. The number $\text{tr} TS$ will also be denoted by

$$(6) \quad T_{\xi_n x_n} B(\xi_1, \dots, \xi_n).$$

The number (6) does not depend on the bound variables ξ_n, x_n but it depends on $\xi_1, \dots, \xi_{n-1}, x_1, \dots, x_{n-1}$, viz. it is a $(2n-2)$ -linear functional $B_1(\xi_1, \dots, \xi_{n-1})$ on $\mathcal{E}^{n-1} \times X^{n-1}$. Suppose that B_1 interpreted as a function of ξ_{n-1} and x_{n-1} only is a Hilbert-Schmidt operator, and that $T' \in \mathfrak{S}$. Then

$$T'_{\xi_{n-1} x_{n-1}} T_{\xi_n x_n} B(\xi_1, \dots, \xi_n)$$

is the number $T'_{\xi_{n-1} x_{n-1}} B_1(\xi_1, \dots, \xi_{n-1})$. Continuing this procedure we can define, under similar hypotheses, the expression

$$T_{r \xi_r x_r} \dots T_{1 \xi_1 x_1} B(\xi_1, \dots, \xi_n) \quad (r \leq n)$$

for $T_1, \dots, T_r \in \mathfrak{S}$. It is not difficult to verify that the required hypotheses are fulfilled in the case where

$$B(\xi_1, \dots, \xi_n) = \begin{vmatrix} 0 & \xi_1 x_2 & \xi_1 x_3 & \dots & \xi_1 x_n \\ \xi_2 x_1 & 0 & \xi_2 x_3 & \dots & \xi_2 x_n \\ \xi_3 x_1 & \xi_3 x_2 & 0 & \dots & \xi_3 x_n \\ \dots & \dots & \dots & \dots & \dots \\ \xi_n x_1 & \xi_n x_2 & \xi_n x_3 & \dots & 0 \end{vmatrix}.$$

In the sequel we shall consider some analytic functions $D(T)$ on \mathfrak{S} and some analytic mappings $D(T)$ from \mathfrak{S} into a Banach space. Then $D'(T; T_1)$ will denote the first differential of $D(T)$, i. e.

$$D'(T; T_1) = \lim_{\varepsilon \rightarrow 0} \varepsilon^{-1} (D(T + \varepsilon T_1) - D(T)).$$

By induction,

$$\begin{aligned} D^{(m)}(T; T_1, \dots, T_m) \\ = \lim_{\varepsilon \rightarrow 0} \varepsilon^{-1} (D(T + \varepsilon T_m; T_1, \dots, T_{m-1}) - D(T; T_1, \dots, T_{m-1})) \end{aligned}$$

for $T, T_1, \dots, T_m \in \mathfrak{S}$. Clearly $D^{(m)}(T; T_1, \dots, T_m)$ is linear and symmetric in variables T_1, \dots, T_m .

If values of $D(T)$ belong to \mathfrak{D}_n , then also $D^{(m)}(T; T_1, \dots, T_m) \in \mathfrak{D}_n$, and $D^{(m)}(T; T_1, \dots, T_m)(\xi_1, \dots, \xi_n)$ is the value of the $2n$ -linear functional $D^{(m)}(T; T_1, \dots, T_m)$ at the point $(\xi_1, \dots, \xi_n, x_1, \dots, x_n) \in \mathcal{E}^n \times X^n$, according to the notation assumed at the beginning of this section. The symbol $D(T)(\xi_1, \dots, \xi_n)$ has the analogous meaning.

§ 2. The Carleman determinant. Let $T \in \mathfrak{D}$. By the Carleman determinant of the operator $A = I + T$ we understand the number

$$(7) \quad D_0(T) = \sum_{m=0}^{\infty} \frac{1}{m!} D_{0,m}(T),$$

where $D_{0,0}(T) = 1$ and, for $m > 0$,

$$(8) \quad D_{0,m}(T) = \begin{vmatrix} 0 & m-1 & 0 & 0 & \dots & 0 & 0 & 0 \\ \text{tr}(T^2) & 0 & m-2 & 0 & \dots & 0 & 0 & 0 \\ \text{tr}(T^3) & \text{tr}(T^2) & 0 & m-3 & \dots & 0 & 0 & 0 \\ \dots & \dots & \dots & \dots & \dots & \dots & \dots & \dots \\ \text{tr}(T^{m-1}) & \text{tr}(T^{m-2}) & \dots & \text{tr}(T^2) & 0 & 1 & & \\ \text{tr}(T^m) & \text{tr}(T^{m-1}) & \dots & \text{tr}(T^3) & \text{tr}(T^2) & 0 & & \end{vmatrix},$$

i. e., by an easy verification,

$$(9) \quad D_{0,m}(T) = T_{\xi_1 x_1} \dots T_{\xi_m x_m} \begin{vmatrix} 0 & \xi_1 x_2 & \xi_1 x_3 & \dots & \xi_1 x_m \\ \xi_2 x_1 & 0 & \xi_2 x_3 & \dots & \xi_2 x_m \\ \xi_3 x_1 & \xi_3 x_2 & 0 & \dots & \xi_3 x_m \\ \dots & \dots & \dots & \dots & \dots \\ \xi_m x_1 & \xi_m x_2 & \xi_m x_3 & \dots & 0 \end{vmatrix}.$$

Carleman [1] and Smithies [1] have proved that

$$(10) \quad |D_{0,m}(T)| \leq m! \left(\frac{e}{m}\right)^{m/2} \|T\|^m.$$

Thus the series (7) converges absolutely. It follows from (9) that $D_{0,m}(T)$ is a homogeneous polynomial of T of the degree m . Hence it follows that $D_0(T)$ is an analytic function defined on the Banach space \mathfrak{S} . Carleman [1] and Smithies [1] have proved that

$$(11) \quad |D_0(T)| \leq \exp(\frac{1}{2}\|T\|^2).$$

By the general theory of homogeneous polynomials (see e. g. Hille and Phillips [1], Chapter XXVI), there exists an m -linear functional $\theta_m(T_1, \dots, T_m)$ on \mathfrak{S}^m , symmetric in variables $T_1, \dots, T_m \in \mathfrak{S}$, such that

$$(12) \quad D_{0,m}(T) = \theta_m(T, \dots, T).$$

The functional θ_m is uniquely determined by $D_{0,m}$, viz.

$$(13) \quad \theta_m(T_1, \dots, T_m) = \frac{1}{m!} \sum (-1)^{i_1 + \dots + i_{m+m}} D_{0,m}(i_1 T_1 + \dots + i_m T_m),$$

where the summation is extended over all sequences i_1, \dots, i_m composed of numbers 0 and 1. It follows easily from (10) and (13) that the norm

$$|\theta_m| = \sup_{\|T_1\| \leq 1, \dots, \|T_m\| \leq 1} |\theta_m(T_1, \dots, T_m)|$$

of θ_m satisfies the inequality

$$(14) \quad |\theta_m| \leq m^{m/2} (2\sqrt{e})^m.$$

We can write an explicit formula for θ_m :

$$(15) \quad \theta_m(T_1, \dots, T_m) = T_{\xi_1 x_1} \dots T_{\xi_m x_m} \begin{vmatrix} 0 & \xi_1 x_2 & \dots & \xi_1 x_m \\ \xi_2 x_1 & 0 & \dots & \xi_2 x_m \\ \dots & \dots & \dots & \dots \\ \xi_m x_1 & \xi_m x_2 & \dots & 0 \end{vmatrix}.$$

In fact, the right side of (15) is a symmetric functional satisfying (12) on account of (9). The detailed proof of the symmetry is similar to an argument in Leżański [1], p. 248.

Since

$$D_{0,m}^{(n)}(T; T_1, \dots, T_n) = \begin{cases} \frac{m!}{(m-n)!} \theta_m(T, \dots, T, T_1, \dots, T_n) & \text{for } n \leq m, \\ 0 & \text{for } n > m, \end{cases}$$

we have

$$(16) \quad D^{(n)}(T; T_1, \dots, T_n) = \sum_{m=0}^{\infty} \frac{1}{m!} \theta_{m+n}(T, \dots, T, T_1, \dots, T_n).$$

By (14),

$$|\theta_{m+n}(T, \dots, T, T_1, \dots, T_n)| \leq (n+m)^{\frac{n+m}{2}} (2\sqrt{e})^{n+m} \|T\|^m \|T_1\| \dots \|T_n\|.$$

Putting $T_1 = x_1 \cdot \xi_1, \dots, T_n = x_n \cdot \xi_n$ we obtain (see (4))

$$(17) \quad |\theta_{m+n}(T, \dots, T, x_1 \cdot \xi_1, \dots, x_n \cdot \xi_n)| \leq (n+m)^{\frac{n+m}{2}} (2\sqrt{e})^{n+m} \|T\|^m |x_1| \cdot |\xi_1| \dots |x_n| \cdot |\xi_n|.$$

For every fixed $T \in \mathfrak{S}$ the expressions $\theta_{m+n}(T, \dots, T, x_1 \cdot \xi_1, \dots, x_n \cdot \xi_n)$ and

$$(18) \quad D_0^{(n)}(T; x_1 \cdot \xi_1, \dots, x_n \cdot \xi_n) = \sum_{m=0}^{\infty} \frac{1}{m!} \theta_{m+n}(T, \dots, T, x_1 \cdot \xi_1, \dots, x_n \cdot \xi_n)$$

(see (16)) are $2n$ -linear bounded functionals on $\mathcal{E}^n \times X^n$, i. e. elements of \mathfrak{D}_n . It follows from (17) that, for every fixed $T \in \mathfrak{S}$, the series of elements of \mathfrak{D}_n on the right side of (18) converges in norm to the left side of (18).

§ 3. The Carleman subdeterminants. By the *first Carleman subdeterminant* of an operator $A = I + T$ ($T \in \mathfrak{S}$) we shall understand the bilinear functional $D_1(T)$ defined by the equality

$$(19) \quad \xi D_1(T)x = D_1(T) \begin{pmatrix} \xi \\ x \end{pmatrix} = D_0'(T; x \cdot \xi) + D_0(T) \cdot \xi x.$$

By induction, the n -th *Carleman subdeterminant* of $A = I + T$ is the $2n$ -linear functional $D_n(T)$ defined by the equality

$$(20) \quad D_n(T) \begin{pmatrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{pmatrix} = D_{n-1}'(T; x_n \cdot \xi_n) \begin{pmatrix} \xi_1, \dots, \xi_{n-1} \\ x_1, \dots, x_{n-1} \end{pmatrix} + D_{n-1}(T) \begin{pmatrix} \xi_1, \dots, \xi_{n-1} \\ x_1, \dots, x_{n-1} \end{pmatrix} \cdot \xi_n x_n.$$

The differentiation on the right side is feasible because one proves, by induction, that $D_n(T)$ is an analytic mapping from \mathfrak{S} into \mathfrak{D}_n . Consequently

$$(21) \quad D_n(T) = \sum_{m=0}^{\infty} \frac{1}{m!} D_{n,m}(T),$$

where $D_{n,m}(T)$ is a homogeneous polynomial (of the variable $T \in \mathfrak{S}$) of the degree m , with values in \mathfrak{D}_n . It is not difficult to prove by induction on n that

$$(22) \quad D_{n,m}(T) \begin{pmatrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{pmatrix} = T_{\xi_{n+1}x_{n+1}} \dots T_{\xi_{n+m}x_{n+m}} \begin{vmatrix} \xi_1 x_1 & \dots & \xi_1 x_n & \xi_1 x_{n+1} & \xi_1 x_{n+2} & \dots & \xi_1 x_{n+m} \\ \dots & \dots & \dots & \dots & \dots & \dots & \dots \\ \xi_n x_1 & \dots & \xi_n x_n & \xi_n x_{n+1} & \xi_n x_{n+2} & \dots & \xi_n x_{n+m} \\ \xi_{n+1} x_1 & \dots & \xi_{n+1} x_n & 0 & \xi_{n+1} x_{n+2} & \dots & \xi_{n+1} x_{n+m} \\ \xi_{n+2} x_1 & \dots & \xi_{n+2} x_n & \xi_{n+2} x_{n+1} & 0 & \dots & \xi_{n+2} x_{n+m} \\ \dots & \dots & \dots & \dots & \dots & \dots & \dots \\ \xi_{n+m} x_1 & \dots & \xi_{n+m} x_n & \xi_{n+m} x_{n+1} & \xi_{n+m} x_{n+2} & \dots & 0 \end{vmatrix}$$

Hence it follows that

$$(23) \quad D_{n,m}(T) = \begin{vmatrix} T_n^0, & m, & 0, & 0, & 0, & \dots, & 0, & 0 \\ T_n^1, & 0, & m-1, & 0, & 0, & \dots, & 0, & 0 \\ T_n^2, & \text{tr } T^2, & 0, & m-2, & 0, & \dots, & 0, & 0 \\ T_n^3, & \text{tr } T^3, & \text{tr } T^2, & 0, & m-3, & \dots, & 0, & 0 \\ \dots & \dots & \dots & \dots & \dots & \dots & \dots & \dots \\ T_n^{m-1}, & \text{tr } T^{m-1}, & \text{tr } T^{m-2}, & \dots, & 0, & 1 \\ T_n^m, & \text{tr } T^m, & \text{tr } T^{m-1}, & \dots, & \text{tr } T^2, & 0 \end{vmatrix}$$

where

$$(24) \quad T_n^m \begin{pmatrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{pmatrix} = \sum \begin{vmatrix} \xi_1 T^{i_1} x_1, \dots, \xi_1 T^{i_1} x_1 \\ \dots \\ \xi_n T^{i_n} x_1, \dots, \xi_n T^{i_n} x_1 \end{vmatrix} = \sum \begin{vmatrix} \xi_1 T^{i_1} x_1, \dots, \xi_1 T^{i_n} x_n \\ \dots \\ \xi_1 T^{i_1} x_1, \dots, \xi_1 T^{i_n} x_n \end{vmatrix}$$

the summation being extended over all sequences i_1, \dots, i_n of non-negative integers whose sum is equal to m .

It follows from (22) or (23) that $D_{n,m}(T) \begin{pmatrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{pmatrix}$ is skew symmetric in ξ_1, \dots, ξ_n and in x_1, \dots, x_n . Thus, by (21), $D_n(T) \begin{pmatrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{pmatrix}$ is also skew symmetric in ξ_1, \dots, ξ_n and in x_1, \dots, x_n .

We can write immediate formulae for $D_n(T)$:

$$(25) \quad D_n(T) \begin{pmatrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{pmatrix} = D_0^{(n)}(T; x_1 \cdot \xi_1, \dots, x_n \cdot \xi_n) + \sum_{i=1}^n D_0^{(n-1)}(T; x_1 \cdot \xi_1, \dots, x_{i-1} \cdot \xi_{i-1}, x_{i+1} \cdot \xi_{i+1}, \dots, x_n \cdot \xi_n) \cdot \xi_i x_i + \sum_{\substack{i,j=1 \\ i < j}}^n D_0^{(n-2)}(T; x_1 \cdot \xi_1, \dots, x_{i-1} \cdot \xi_{i-1}, x_{i+1} \cdot \xi_{i+1}, \dots, x_{j-1} \cdot \xi_{j-1}, x_{j+1} \cdot \xi_{j+1}, \dots, x_n \cdot \xi_n) \cdot \xi_i x_i \cdot \xi_j x_j + \dots + \sum_{i=1}^n D_0'(T; x_i \cdot \xi_i) \cdot \xi_1 x_1 \cdot \dots \cdot \xi_{i-1} x_{i-1} \cdot \xi_{i+1} x_{i+1} \cdot \dots \cdot \xi_n x_n + D_0(T) \cdot \xi_1 x_1 \cdot \dots \cdot \xi_n x_n.$$

It follows from (18), (21) and (25) that

$$(26) \quad D_{n,m}(T) \begin{pmatrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{pmatrix} = \theta_{m+n}(T, \dots, T, x_1 \cdot \xi_1, \dots, x_n \cdot \xi_n) + \sum_{i=1}^n \theta_{m+n-1}(T, \dots, T, x_1 \cdot \xi_1, \dots, x_{i-1} \cdot \xi_{i-1}, x_{i+1} \cdot \xi_{i+1}, \dots, x_n \cdot \xi_n) \cdot \xi_i x_i + \sum_{\substack{i,j=1 \\ i < j}}^n \theta_{m+n-2}(T, \dots, T, x_1 \cdot \xi_1, \dots, x_{i-1} \cdot \xi_{i-1}, x_{i+1} \cdot \xi_{i+1}, \dots, x_{j-1} \cdot \xi_{j-1}, x_{j+1} \cdot \xi_{j+1}, \dots, x_n \cdot \xi_n) \cdot \xi_i x_i \cdot \xi_j x_j + \dots + \sum_{i=1}^n \theta_{m+1}(T, \dots, T, x_i \cdot \xi_i) \cdot \xi_1 x_1 \cdot \dots \cdot \xi_{i-1} x_{i-1} \cdot \xi_{i+1} x_{i+1} \cdot \dots \cdot \xi_n x_n + \theta_m(T, \dots, T) \cdot \xi_1 x_1 \cdot \dots \cdot \xi_n x_n.$$

By (17), and (26) we have the following estimation of the norm of the $2n$ -linear functional $D_{n,m}(T)$:

$$(27) \quad |D_{n,m}(T)| \leq 2^n (n+m)^{\frac{n+m}{2}} (2\sqrt{e})^{n+m} \|T\|^m.$$

This proves that, for every fixed $T \in \mathfrak{S}$, the series (21) on elements of \mathfrak{D}_n converges in norm.

Formula (25) expresses $D_n(T)$ by means of $D_0(T)$ and its derivatives. The following formula expresses the derivatives of $D_0(T)$ by means of D_0, \dots, D_n, \dots :

$$\begin{aligned}
 (28) \quad D_0^{(n)}(T; x_1 \cdot \xi_1, \dots, x_n \cdot \xi_n) &= D_n(T) \left(\begin{matrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{matrix} \right) - \\
 &- \sum_{i=1}^n D_{n-1}(T) \left(\begin{matrix} \xi_1, \dots, \xi_{i-1}, \xi_{i+1}, \dots, \xi_n \\ x_1, \dots, x_{i-1}, x_{i+1}, \dots, x_n \end{matrix} \right) \cdot \xi_i x_i + \\
 &+ \sum_{\substack{i,j=1 \\ i < j}}^n D_{n-2}(T) \left(\begin{matrix} \xi_1, \dots, \xi_{i-1}, \xi_{i+1}, \dots, \xi_{j-1}, \xi_{j+1}, \dots, \xi_n \\ x_1, \dots, x_{i-1}, x_{i+1}, \dots, x_{j-1}, x_{j+1}, \dots, x_n \end{matrix} \right) \cdot \xi_i x_i \cdot \xi_j x_j - \\
 &- \dots \dots \dots + \\
 &+ (-1)^{n-1} D_1(T) \left(\begin{matrix} \xi_i \\ x_i \end{matrix} \right) \cdot \xi_1 x_1 \cdot \dots \cdot \xi_{i-1} x_{i-1} \cdot \xi_{i+1} x_{i+1} \cdot \dots \cdot \xi_n x_n + \\
 &+ (-1)^n D_0(T) \cdot \xi_1 x_1 \cdot \dots \cdot \xi_n x_n.
 \end{aligned}$$

Formula (28) follows immediately from (20) and (19).

§ 4. Applications to the theory of linear equations. First we shall prove that

(i) For every Hilbert-Schmidt operator T , the sequence

$$(29) \quad D_0(T), D_1(T), D_2(T), \dots$$

is a determinant system for the operator $A = I + T$.

For the definition of determinant system—see Sikorski [3] p. 172, conditions (d₁)-(d₅). We have observed in § 3 that (d₁) and (d₃) hold. (d₃) also holds because X, \mathcal{E} are Hilbert spaces.

We have to prove (d₄), i. e. that, for every fixed $T \in \mathfrak{S}$, at least one of the multilinear functionals (29) does not vanish identically. Suppose the contrary, i. e. that, for a $T \in \mathfrak{S}$, we have $D_n(T) = 0$ for $n = 0, 1, 2, \dots$. Hence, by (28),

$$(30) \quad D_0^{(n)}(T; T_1, \dots, T_n) = 0$$

for $n = 0, 1, 2, \dots$ and for all one-dimensional operators T_1, \dots, T_n . Since $D_0^{(n)}(T; T_1, \dots, T_n)$ is linear in T_1, \dots, T_n , (30) holds for all finitely dimensional operators T_1, \dots, T_n . Since $D_0^{(n)}(T; T_1, \dots, T_n)$ is continuous in $T_1, \dots, T_n \in \mathfrak{S}$ and finitely dimensional operators are dense in \mathfrak{S} , (30) holds for every n and for all $T_1, T_2, \dots, T_n \in \mathfrak{S}$. Since the analytic function $D_0(T)$ vanishes at the point T with all its derivatives, it is equal to zero everywhere. This is impossible since $D_0(0) = 1$.

Now we have to prove (d₅), i. e. the identities

$$\begin{aligned}
 (31) \quad D_{n+1}(T) \left(\begin{matrix} \xi_0 A, \xi_1, \dots, \xi_n \\ x_0, x_1, \dots, x_n \end{matrix} \right) \\
 = \sum_{i=0}^n (-1)^i \xi_0 x_i \cdot D_n(T) \left(\begin{matrix} \xi_0, \dots, \xi_n \\ x_0, \dots, x_{i-1}, x_{i+1}, \dots, x_n \end{matrix} \right),
 \end{aligned}$$

$$\begin{aligned}
 (32) \quad D_{n+1}(T) \left(\begin{matrix} \xi_0, \xi_1, \dots, \xi_n \\ Ax_0, x_1, \dots, x_n \end{matrix} \right) \\
 = \sum_{i=0}^n (-1)^i \xi_i x_0 \cdot D_n(T) \left(\begin{matrix} \xi_0, \dots, \xi_{i-1}, \xi_{i+1}, \dots, \xi_n \\ x_1, \dots, x_n \end{matrix} \right).
 \end{aligned}$$

By the same argument as in Leżański [1], p. 254-256, we prove first those identities for $D_{n+1,m}(T)$, $D_{n,m}(T)$ instead of $D_{n+1}(T)$, $D_n(T)$ respectively. Hence we obtain (31), (32).

The next theorem follows immediately from (i) and Sikorski [3], Theorem II.

(ii) Let T be a Hilbert-Schmidt operator, and let r be the smallest integer such that the $2r$ -linear functional D_r does not vanish identically. Let $\eta_1, \dots, \eta_r, y_1, \dots, y_r$ be such that

$$D_r(T) \left(\begin{matrix} \eta_1, \dots, \eta_r \\ y_1, \dots, y_r \end{matrix} \right) \neq 0.$$

Then there exist elements $\xi_1, \dots, \xi_r, z_1, \dots, z_r$ and an operator B such that, for all ξ, x

$$\xi_i x = \frac{D_r(T) \left(\begin{matrix} \eta_1, \dots, \eta_r \\ y_1, \dots, y_{i-1}, x, y_{i+1}, \dots, y_r \end{matrix} \right)}{D_r(T) \left(\begin{matrix} \eta_1, \dots, \eta_r \\ y_1, \dots, y_r \end{matrix} \right)},$$

$$\xi z_i = \frac{D_r(T) \left(\begin{matrix} \eta_1, \dots, \eta_{i-1}, \xi, \eta_{i+1}, \dots, \eta_r \\ y_1, \dots, y_r \end{matrix} \right)}{D_r(T) \left(\begin{matrix} \eta_1, \dots, \eta_r \\ y_1, \dots, y_r \end{matrix} \right)},$$

$$\xi Bx = \frac{D_{r+1}(T) \left(\begin{matrix} \xi, \eta_1, \dots, \eta_r \\ x, y_1, \dots, y_r \end{matrix} \right)}{D_r(T) \left(\begin{matrix} \eta_1, \dots, \eta_r \\ y_1, \dots, y_r \end{matrix} \right)}.$$

The elements ξ_1, \dots, ξ_r are linearly independent in \mathcal{E} , and so are z_1, \dots, z_r in X .

by the same argument as in the case of (39).

$$(53) \quad D_1^*(T) = TD_1(T) = D_1(T)T,$$

$D_1^*(T)$ is a Hilbert-Schmidt operator, and the series

$$(54) \quad D_1^*(T) = \sum_{m=0}^{\infty} D_{1,m}^*(T)$$

of Hilbert-Schmidt operators converges in the norm $\|\cdot\|$ in \mathfrak{S} .

By (52), (48) and (21)

$$(55) \quad D_n^*(T) \begin{pmatrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{pmatrix} = D_n(T) \begin{pmatrix} \xi_1 T, \dots, \xi_n T \\ x_1, \dots, x_n \end{pmatrix} = D_n(T) \begin{pmatrix} \xi_1, \dots, \xi_n \\ T x_1, \dots, T x_n \end{pmatrix}.$$

Let

$$D_0^*(T) = D_0(T).$$

The sequence

$$D_0^*(T), D_1^*(T), D_2^*(T), \dots$$

is not a determinant system for the operator $A = I + T$ in the sense defined by Sikorski [3] since, instead of (31), (32), it satisfies the following identities:

$$(56) \quad D_{n+1}^*(T) \begin{pmatrix} \xi_0 A, \xi_1, \dots, \xi_n \\ x_0, x_1, \dots, x_n \end{pmatrix} = \sum_{i=0}^n (-1)^i \xi_0 T x_i \cdot D_n^*(T) \begin{pmatrix} \xi_0, \dots, \xi_n \\ x_0, \dots, x_{i-1}, x_{i+1}, \dots, x_n \end{pmatrix},$$

$$(57) \quad D_{n+1}^*(T) \begin{pmatrix} \xi_0, \xi_1, \dots, \xi_n \\ A x_0, x_1, \dots, x_n \end{pmatrix} = \sum_{i=0}^n (-1)^i \xi_i T x_0 \cdot D_n^*(T) \begin{pmatrix} \xi_0, \dots, \xi_{i-1}, \xi_{i+1}, \dots, \xi_n \\ x_1, \dots, x_n \end{pmatrix}.$$

However, it can be used to solve the equations $Ax = x_0$ and $\xi A = \xi_0$. By the same method as in Sikorski [4] § 3 we can prove the following theorem:

(vii) Let T be a Hilbert-Schmidt operator. The smallest integer r such that $D_r^*(T) \neq 0 \in \mathfrak{D}_r$ is equal to the smallest integer r such that $D_r(T) \neq 0 \in \mathfrak{D}_r$. Then $D_r^*(T) = (-1)^r D_r(T)$. Let $\eta_1, \dots, \eta_r, y_1, \dots, y_r$ be such that $D_r^* \begin{pmatrix} \eta_1, \dots, \eta_r \\ y_1, \dots, y_r \end{pmatrix} \neq 0$. Then there exist elements $\xi_1^*, \dots, \xi_r^*, z_1^*, \dots, z_r^*$ and a Hilbert-Schmidt operator B^* such that for all ξ, x

$$\xi_i^* x = \frac{D_r^*(T) \begin{pmatrix} \eta_1, \dots, \eta_r \\ y_1, \dots, y_{i-1}, x, y_{i+1}, \dots, y_r \end{pmatrix}}{D_r^*(T) \begin{pmatrix} \eta_1, \dots, \eta_r \\ y_1, \dots, y_r \end{pmatrix}},$$

$$\xi z_i^* = \frac{D_r^*(T) \begin{pmatrix} \eta_1, \dots, \eta_{i-1}, \xi, \eta_{i+1}, \dots, \eta_r \\ y_1, \dots, y_r \end{pmatrix}}{D_r^*(T) \begin{pmatrix} \eta_1, \dots, \eta_r \\ y_1, \dots, y_r \end{pmatrix}},$$

$$\xi B^* x = \frac{D_{r+1}^*(T) \begin{pmatrix} \xi, \eta_1, \dots, \eta_r \\ x, y_1, \dots, y_r \end{pmatrix}}{D_r^*(T) \begin{pmatrix} \eta_1, \dots, \eta_r \\ y_1, \dots, y_r \end{pmatrix}}.$$

The elements ξ_1^*, \dots, ξ_r^* are linearly independent in \mathfrak{E} , and so are z_1^*, \dots, z_r^* in X .

The equation

$$x + Tx = x_0$$

has a solution x if and only if $\xi_i^* x_0 = 0$ for $i = 1, \dots, r$. Then the general form of the solution is

$$x = x_0 - B^* x_0 + c_1 z_1^* + \dots + c_r z_r^*.$$

The adjoint equation

$$\xi + \xi T = \xi_0$$

has a solution ξ if and only if $\xi_0 z_i^* = 0$ for $i = 1, \dots, r$. Then the general form of the solution is

$$\xi = \xi_0 - \xi_0 B^* + c_1 \xi_1^* + \dots + c_r \xi_r^*.$$

In the case of $r = 0$ Theorem (vii) asserts that

$$(58) \quad I - \frac{D_1^*(T)}{D_0^*(T)} = (I + T)^{-1}.$$

§ 7. An integral model. Let μ be a measure defined on a σ -field of subsets of a set Γ . The integrals taken over the whole space Γ will be denoted, for brevity, by $\int f(t) dt$ instead of $\int f(t) d\mu(t)$, and similarly for multiple integrals.

Suppose now that $X = \mathfrak{E} = L^2(\Gamma, \mu)$. The class \mathfrak{S} of all Hilbert-Schmidt operators coincides (see e.g. Smithies [1]) with the class of all integral operators T

$$Tx(s) = \int \tau(s, t) x(t) dt$$

where

$$(59) \quad \|T\| = \sqrt{\iint |\tau(s, t)|^2 ds dt} < \infty,$$

i. e. $\tau \in L^2(\Gamma \times \Gamma, \mu \times \mu)$.

Similarly as in Sikorski [4], § 4, we can prove that

$$(60) \quad D_{n,m}^*(T) \left(\begin{matrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{matrix} \right) \\ = \int \dots \int \vartheta_{n,m}^* \left(\begin{matrix} s_1, \dots, s_n \\ t_1, \dots, t_n \end{matrix} \right) \xi_1(s_1) \dots \xi_n(s_n) \cdot x_1(t_1) \dots x_n(t_n) ds_1 \dots ds_n dt_1 \dots dt_n,$$

where

$$(61) \quad \vartheta_{n,m}^* \begin{pmatrix} s_1, \dots, s_n \\ t_1, \dots, t_n \end{pmatrix} = \int \dots \int \begin{vmatrix} T(s_1, t_1), \dots, T(s_1, t_n), & T(s_1, r_1), & T(s_1, r_2)T, & (s_1, r_3), \dots, & T(s_1, r_m) \\ \vdots & \vdots & \vdots & \vdots & \vdots \\ T(s_n, t_1), \dots, T(s_n, t_n), & T(s_n, r_1), & T(s_n, r_2), & T(s_n, r_3), \dots, & T(s_n, r_m) \\ \vdots & \vdots & \vdots & \vdots & \vdots \\ T(r_1, t_1), \dots, T(r_1, t_n), & 0, & T(r_1, r_2), & T(r_1, r_3), \dots, & T(r_1, r_m) \\ \vdots & \vdots & \vdots & \vdots & \vdots \\ T(r_2, t_1), \dots, T(r_2, t_n), & T(r_2, r_1), & 0, & T(r_2, r_3), \dots, & T(r_2, r_m) \\ \vdots & \vdots & \vdots & \vdots & \vdots \\ T(r_3, t_1), \dots, T(r_3, t_n), & T(r_3, r_1), & T(r_3, r_2), & 0, & \dots, & T(r_3, r_m) \\ \vdots & \vdots & \vdots & \vdots & \vdots & \vdots \\ T(r_m, t_1), \dots, T(r_m, t_n), & T(r_m, r_1), & T(r_m, r_2), & T(r_m, r_3), \dots, & 0 \end{vmatrix} \times dr_1 \dots dr_m.$$

Thus formally

$$\begin{aligned}
(62) \quad & D_n^*(T) \left(\begin{matrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{matrix} \right) \\
&= \int \dots \int \partial_n^* \left(\begin{matrix} s_1, \dots, s_n \\ t_1, \dots, t_n \end{matrix} \right) \xi_1(s_1) \dots \xi_n(s_n) x_1(t_1) \dots x_n(t_n) ds_1 \dots ds_n dt_1 \dots dt_n
\end{aligned}$$

where

$$(63) \quad \vartheta_n^* \left(\begin{matrix} s_1, \dots, s_n \\ t_1, \dots, t_n \end{matrix} \right) = \sum_{m=0}^{\infty} \frac{1}{m!} \vartheta_{n,m}^* \left(\begin{matrix} s_1, \dots, s_n \\ t_1, \dots, t_n \end{matrix} \right).$$

The problem under which conditions the series (63) converges, in some sense, to a function $\vartheta_{t_1, \dots, t_n}^*(s_1, \dots, s_n)$ and relation (62) holds, has not been investigated in detail. For $n = 1$, it follows from § 6 (54) that $\vartheta_{1,m}^* \in L^2(\Gamma \times \Gamma, \mu \times \mu)$, the series

$$\vartheta_1^* \begin{pmatrix} s \\ t \end{pmatrix} = \sum_{m=0}^{\infty} \vartheta_{1,m}^* \begin{pmatrix} s \\ t \end{pmatrix}$$

converges in $L^2(\Gamma \times \Gamma, \mu \times \mu)$ and the function $\vartheta_1^* \in L^2(\Gamma \times \Gamma, \mu \times \mu)$ satisfies the identity:

$$D_1^*(T) \begin{pmatrix} \xi \\ \eta \end{pmatrix} = \iint \vartheta_1^* \begin{pmatrix} s \\ t \end{pmatrix} \xi(s) \eta(t) ds dt.$$

If every one-point subset of Γ has a positive measure, then the series (63) converges pointwise and the identity (62) holds ($n = 1, 2, \dots$).

To give an integral formula for $D_n(T)$ we have to introduce, similarly as in Sikorski [4], § 4, a substitute of the Dirac delta distribution, i. e. a formal expression $\delta(s, t)$ which is a formal kernel of the operator I . By definition,

$$\int \delta(s, t)x(t)dt = x(s), \quad \int \xi(s)\delta(s, t)ds = \xi(t),$$

$$\iiint \delta(s, t) \tau_1(t, r) \tau_2(r, s) dr ds dt = \iint \tau_1(t, s) \tau(s, t) ds dt,$$

etc. (for details, see Sikorski [4], § 4). Then we can write formally

$$(64) \quad D_{n,m}(T) \left(\begin{matrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{matrix} \right) \\ = \int \dots \int \vartheta_{n,m} \left(\begin{matrix} s_1, \dots, s_n \\ t_1, \dots, t_n \end{matrix} \right) \xi_1(s_1) \dots \xi_n(s_n) x_1(t_1) \dots x_n(t_n) ds_1 \dots ds_n dt_1 \dots dt_n,$$

where

$$(65) \quad \partial_{n,m} \begin{pmatrix} s_1, \dots, s_n \\ t_1, \dots, t_n \end{pmatrix} = \int \dots \int \left| \begin{array}{ccccccccc} \delta(s_1, t_1), \dots, \delta(s_1, t_n), \delta(s_1, r_1), \delta(s_1, r_2), \delta(s_1, r_3), \dots, \delta(s_1, r_m) \\ \vdots \\ \delta(s_n, t_1), \dots, \delta(s_n, t_n), \delta(s_n, r_1), \delta(s_n, r_2), \delta(s_n, r_3), \dots, \delta(s_n, r_m) \\ \tau(r_1, t_1), \dots, \tau(r_1, t_n), & 0, & \tau(r_1, r_2), \tau(r_1, r_3), \dots, \tau(r_1, r_m) \\ \tau(r_2, t_1), \dots, \tau(r_2, t_n), \tau(r_2, r_1), & 0, & \tau(r_2, r_3), \dots, \tau(r_2, r_m) \\ \tau(r_3, t_1), \dots, \tau(r_3, t_n), \tau(r_3, r_1), \tau(r_3, r_2), & 0, & \dots, \tau(r_3, r_m) \\ \vdots \\ \tau(r_m, t_1), \dots, \tau(r_m, t_n), \tau(r_m, r_1), \tau(r_m, r_2), \tau(r_m, r_3), \dots, & 0 \end{array} \right| \times dr_1 \dots dr_m$$

and

$$(66) \quad D_n(T) \left(\begin{matrix} \xi_1, \dots, \xi_n \\ x_1, \dots, x_n \end{matrix} \right) \\ = \int \dots \int \vartheta_n \left(\begin{matrix} s_1, \dots, s_n \\ t_1, \dots, t_n \end{matrix} \right) \xi_1(s_1) \dots \xi_n(s_n) x_1(t_1) \dots x_n(t_n) ds_1 \dots ds_n dt_1 \dots dt_n$$

where

$$(67) \quad \vartheta_n \left(\begin{smallmatrix} s_1, \dots, s_n \\ t_1, \dots, t_n \end{smallmatrix} \right) = \sum_{m=0}^{\infty} \frac{1}{m!} \vartheta_{n,m} \left(\begin{smallmatrix} s_1, \dots, s_n \\ t_1, \dots, t_n \end{smallmatrix} \right).$$

If every one-point subset of Γ has a positive measure, then $\delta(s, t)$ and consequently, also $\vartheta_{n,m} \left(\begin{smallmatrix} s_1, \dots, s_n \\ t_1, \dots, t_n \end{smallmatrix} \right)$, $\vartheta_n \left(\begin{smallmatrix} s_1, \dots, s_n \\ t_1, \dots, t_n \end{smallmatrix} \right)$ are functions, and the series (66) converges pointwise ($n = 1, 2, \dots$).

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Réfractions non-hilbertiennes d'une transformation symétrique bornée

par

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1. Introduction. Dans tout ce qui suit, μ est une mesure sur un clan de parties d'un ensemble ω , et L^0 est l'ensemble des fonctions étagées μ -mesurables sur ω . Remarquons que L^0 est partout dense dans l'espace de Banach $\mathfrak{X}_r = \{x: \|x\|_{1/r} < \infty\}$, où

$$\|x\|_{1/r} = \left(\int_{\omega} |x|^{1/r} d\mu \right)^r;$$

nous supposons désormais que $0 < r < 1$. Soit \mathcal{H} l'ensemble des transformations symétriques bornées de l'espace hilbertien $\mathfrak{X}_{1/2}$. Posons $H \in \mathcal{H}$ et

$$|H|_r = \sup \{ \|Hx\|_{1/r} : x \in L^0 \text{ et } \|x\|_{1/r} \leq 1 \}.$$

Soit H^0 la restriction de H à L^0 , et désignons par H_r le prolongement linéaire continu de H^0 à \mathfrak{X}_r ; nous dirons que H_r est la „réfraction” de H . Précisons: H_r est l'unique prolongement de H^0 appartenant à l'ensemble \mathcal{C}_r des endomorphismes continus de \mathfrak{X}_r . Il est clair que H possèdera une réfraction si et seulement si $|H|_r \neq \infty$.

Rappelons qu'il existe une bijection $H \rightarrow E^H$ de \mathcal{H} sur l'ensemble des familles spectrales (cf. [13], pp. 174-176). Nous noterons \mathcal{H}_i l'ensemble des transformations H appartenant à \mathcal{H} telles que

$$(i) \quad \infty \neq \sup_{0 < \lambda < \infty} |E^H(\lambda)|_s \quad \text{lorsque} \quad 0 < s < 1.$$

Nous supposons désormais que $H \in \mathcal{H}_i$. On verra (au n°. 2.8) que H possède une réfraction H_r ; en outre, H_r admet une décomposition spectrale (bien que H_r ne soit pas nécessairement un „spectral operator” dans le sens de N. Dunford: voir n°. 3.4 et § 5). La réfraction hérite de H d'autres propriétés spectrales⁽¹⁾. Par exemple, H_r possède un „calcul

⁽¹⁾ Il serait intéressant de déterminer quelles autres propriétés spectrales de H sont héritées par H_r . A ce propos, voir § 5.