

Positivity and Monotonicity Properties
of Transport Equations with Spatially Dependent
Cross Sections^{*}

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Abstract

We investigate the transport equation

$$\mu \frac{\partial \psi}{\partial x}(x, \mu) + \psi(x, \mu) = \frac{1}{2} \xi c(x) \int_{-1}^{+1} \psi(x, \mu') d\mu' + f(x, \mu) \quad (0 < x < \tau < \infty)$$

with suitable boundary conditions through an equivalent integral equation. Assuming the incoming fluxes, the internal source term $f(x, \mu)$, the cross section $c(x)$ and the parameter ξ to be nonnegative, we prove the existence of a unique dominant eigenvalue $\xi = \xi_0(\tau)$ for which the homogeneous problem has a positive solution (critical case), the existence of a unique positive solution for $\xi < \xi_0(\tau)$ (non-critical case), and the absence of positive solutions for $\xi > \xi_0(\tau)$ (supercritical case). We show $\xi_0(\tau)$ to decrease continuously from ∞ to some $\xi_0(\infty) > 0$ whenever τ increases from 0 to ∞ (monotonicity). The results are obtained by studying an operator that leaves invariant the cone of nonnegative functions in $L_\infty(0, \tau)$.

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Introduction

The time-independent neutron transport equation in a homogeneous finite or semi-infinite slab medium with spatially dependent collision ratio $\xi c(x)$ reads as follows:

$$\mu \frac{\partial \psi}{\partial x}(x, \mu) + \psi(x, \mu) = \frac{1}{2} \xi c(x) \int_{-1}^{+1} \psi(x, \mu') d\mu' + f(x, \mu) \quad (0 < x < \tau). \quad (1)$$

Studied on the finite slab $0 < x < \tau < \infty$ one imposes the boundary conditions

$$\psi(0, \mu) = \phi(\mu) \quad (0 \leq \mu \leq 1), \quad \psi(\tau, \mu) = \phi(\mu) \quad (-1 \leq \mu < 0); \quad (2)$$

on the half-line $0 < x < \tau = \infty$ one imposes the conditions

$$\psi(0, \mu) = \phi(\mu) \quad (0 \leq \mu \leq 1), \quad \int_{-1}^{+1} |\psi(x, \mu)|^2 d\mu = 0(1) \quad (x \rightarrow \infty). \quad (3)$$

The popular method to study Eq. (1) is the method of singular eigenfunction expansion for which the present state of affairs has been reviewed by Larsen.¹ In the present article a different method is adopted that allows one to investigate the existence and uniqueness of nonnegative solutions ψ of Eqs. (1)-(2) and Eqs. (1)-(3), where we assume that $\xi > 0$ and $\phi(\mu)$, $c(x)$ and $f(x, \mu)$ are suitable nonnegative functions. These positivity requirements all are prescribed by the physics behind the mathematical problem.

Let us introduce some notation. By $H = L_2[-1, 1]$ we denote the Hilbert space of square integrable real-valued functions on $[-1, 1]$ with inner product $\langle f, g \rangle = \int_{-1}^{+1} f(\mu)g(\mu)d\mu$. Further, let $\psi(x)$ and $f(x)$ be vectors in H , T and B the self-adjoint operators and Q_+ and Q_- the complementary orthogonal projections defined by

$$\psi(x)(\mu) = \psi(x, \mu), \quad f(x)(\mu) = f(x, \mu); \quad (4)$$

$$(Th)(\mu) = \mu h(\mu), \quad (Bh)(\mu) = \frac{1}{2} \int_{-1}^{+1} h(\mu') d\mu'; \quad (4b)$$

$$(Q_+h)(\mu) = h(\mu) \quad (\mu \geq 0), \quad (Q_-h)(\mu) = 0 \quad (\mu \leq 0). \quad (4c)$$

Then the ranges of Q_{\pm} are the subspaces $\text{Ran } Q_{+} = L_2[0,1]$ and $\text{Ran } Q_{-} = L_2[-1,0]$. The integro-differential equation (1) can now be written in the concise form

$$(T\psi)'(x) = -\psi(x) + \xi c(x)B\psi(x) + f(x) \quad (0 < x < \tau) \quad (5a)$$

with boundary conditions

$$\lim_{x \downarrow 0} \|Q_{+}\psi(x) - Q_{+}\phi\| = \lim_{x \uparrow \tau} \|Q_{-}\psi(x) - Q_{-}\phi\| = 0 \quad (5b)$$

for finite τ with $\phi \in H$, and with boundary conditions

$$\lim_{x \downarrow 0} \|Q_{+}\psi(x) - \phi\| = 0, \quad \|\psi(x)\| = o(1) \quad (x \rightarrow \infty) \quad (5c)$$

for $\tau = \infty$ and $\phi \in \text{Ran } Q_{+} \subset H$. By a solution of the boundary value problem (5a)-(5b) (resp. (5a)-(5c)) we mean a function ψ from $(0, \tau)$ (resp. $(0, \infty)$) into H such that $T\psi$ is strongly differentiable on $(0, \tau)$ (resp. $(0, \infty)$) and satisfies (5a)-(5b) (resp. (5a)-(5c)).

In this article we are searching for nonnegative solutions of (5a)-(5b) and (5a)-(5c). Throughout the paper $c(x)$ will be a bounded measurable nonnegative function on $(0, \tau)$ and K will denote the subset of functions $h \in H = L_2[-1,1]$ such that $h(\mu) > 0$ for almost every $-1 < \mu < 1$. We assume throughout that $\xi > 0$, $\phi \in K$ and f is a bounded uniformly Hölder continuous function² from $[0, \tau]$ (τ finite) or $[0, \infty)$ ($\tau = \infty$) into K . We shall prove the following results:

Theorem 1. Let τ be finite and $c(x) > 0$, except for a set in $(0, \tau)$ of measure zero. Then one of the following three situations occurs:

- (a) non-critical case. For every $\phi \in H$ and bounded uniformly Hölder continuous $f: (0, \tau) \rightarrow H$ there exists a unique solution $\psi: (0, \tau) \rightarrow H$. Whenever $\phi \in K$ and $f(x) \in K$ for $0 < x < \tau$, also the solution ψ assumes its values in K .
- b. critical case. For $\phi = 0$ and $f \equiv 0$ there exists a nontrivial solution $\psi: (0, \tau) \rightarrow K$, which is unique up to multiplication by a positive constant. Whenever $\phi \in K$ and $f(x) \in K$ for $0 < x < \tau$, and at least one of them is non-trivial, there does not exist any solution ψ taking its values in K .

- c. supercritical case. For $\phi \in K$ and $f(x) \in K$ for $0 < x < \tau$, the existence of a solution $\psi: (0, \tau) \rightarrow K$ implies that $\phi = 0$, $f \equiv 0$ and $\psi \equiv 0$. Moreover, there exists a unique positive number $\xi_0(\tau)$, such that for $0 < \xi < \xi_0(\tau)$, $\xi = \xi_0(\tau)$ and $\xi > \xi_0(\tau)$ the non-critical, critical and supercritical cases occur, respectively. The function $\tau \rightarrow \xi_0(\tau)$ is strictly monotonic and continuous with $\lim_{\tau \downarrow 0} \xi_0(\tau) = \infty$ and $\lim_{\tau \rightarrow \infty} \xi_0(\tau) \stackrel{\text{def.}}{=} \xi_0(\infty) > 0$.

Theorem 1 concerns the criticality problem of finding for fixed τ the unique parameter ξ for which (5a)-(5b) with $\phi = 0$ and $f(x) \equiv 0$ has a positive solution. The monotonicity of $\tau \mapsto \xi_0(\tau)$ implies that for every $\xi > \xi_0(\infty)$ there also is a unique τ for which this problem has a positive solution. The former problem one could call the dominant eigenvalue problem; the latter is the actual critical slab problem.

For $c_0(x) \equiv 1$ one has $\xi_0(\infty) = 1$, which corresponds to the $c=1$ half-space problem.³ For this case Busoni, Mangiarotti and Frosali⁴ proved Theorem 1 partially.⁵ In fact, they settled the existence of $\xi_0(\tau)$, the occurrence of the case (a) for $0 < \xi < \xi_0(\tau)$ and the uniqueness of a solution $\psi: (0, \tau) \rightarrow K$ for $\xi = \xi_0(\tau)$, $\phi = 0$ and $f \equiv 0$. In the inhomogeneous case and below criticality Busbridge⁶ studied an integral form of Eq.(1). Kelley⁶ proved part (b) and the existence of $\xi_0(\tau)$ in a different way for positive $c \in L_1(0, \tau)$. Teczan⁷ recently showed interest in the monotonicity of the critical slab value τ_0 for the homogeneous case with a backward scattering term.

The critical half-space problem was dealt with for exponential media ($c(x) = e^{-x/s}$) by Pomraning,⁸ who obtained numerical results. In the present article we prove analogues of Theorem 1, (a) and (c).

Theorem 2. Let $\tau = \infty$ and $c(x) > 0$, except on a set of measure zero. Then one of the following situations occurs:

- (a) $0 < \xi < \xi_0(\infty)$: non-critical case. For every $\phi \in H$ and bounded uniformly Hölder continuous² function $f: (0, \infty) \rightarrow H$ there exists a unique solution $\psi: (0, \infty) \rightarrow H$. Whenever $\phi \in K$ and $f(x) \in K$ for $0 < x < \infty$, the solution ψ also assumes its values in K .
- (c) $\xi > \xi_0(\infty)$: supercritical case. For $\phi \in K$ and $f(x) \in K$ there does not exist a solution $\psi: (0, \infty) \rightarrow K$, unless $\phi = 0$, $f \equiv 0$ and $\psi \equiv 0$.

Both theorems will be proved by reducing the boundary value problems (5a)-(5b) and (5a)-(5c) to a scalar integral equation, which can be written in the form

$$(I - \xi R)\zeta = \bar{\omega}, \tag{6}$$

on the Banach space $BC(0, \tau)$ of bounded continuous functions $\zeta : (0, \tau) \rightarrow \mathbb{R}$ with norm $\|\zeta\| = \sup \{|\zeta(x)| : 0 < x < \tau\}$. In this space we consider the cone (see Section 2 for the definition)

$$K_c = \{\zeta \in BC(0, \tau) : \zeta(x) \geq 0 \text{ for } 0 < x < \tau\}.$$

For $\phi \in K$ and $f(x) \in K$ ($0 < x < \tau$) it appears that $\bar{\omega} \in K_c$ and $R[K_c] \subset K_c$. We are searching for solutions ζ in K_c .

In order to investigate Eq. (6) we shall apply the theory of operators leaving invariant a cone in a Banach space. This theory made considerable progress by the pioneering work of Krein and Rutman,⁹ and Krasnoselskii.¹⁰ Theorems 1 and 2 will appear as consequences of this theory.

In Section 1 the problems (5a)-(5b) and (5a)-(5c) are proved to be equivalent to a scalar integral equation, which will be put into the form (6). Section 2 reviews some of the theory of cones in a Banach space. In Section 3 the proofs of Theorems 1 and 2 will be given, and at the end of the section we prove a statement of Tezcan⁷ rigorously.

1. Statement as a vector equation

Consider the vector $e(\mu) \equiv 1$ in H and define the propagator function $H(x)$ ($0 \neq x \in \mathbb{R}$) by

$$(H(x)h)(\mu) = \begin{cases} |\mu|^{-1} e^{-x/\mu} h(\mu), & x\mu > 0; \\ 0, & x\mu < 0. \end{cases}$$

One easily computes that

$$Ei(x) = \langle H(x)e, e \rangle = \int_1^\infty z^{-1} e^{-z|x|} dz, \tag{7} \quad 0 \neq x \in \mathbb{R}.$$

LEMMA 1. For finite (resp. infinite) τ let f be a bounded

uniformly Hölder continuous² function from $[0, \tau]$ (resp. $[0, \infty)$) into H .

Then

$$F(x) = [I - e^{-xT} Q_+ e^{(\tau-x)T} Q_-] f(x) + \int_0^\tau H(x-y) \cdot \{f(y) - f(x)\} dy,$$

with the second semigroup taken zero for $\tau = \infty$, is a bounded continuous function from $[0, \tau]$ (resp. $[0, \infty)$) into H . Further, $TF(x)$ is strongly differentiable on $(0, \tau)$ and

$$(TF)'(x) = -F(x) + f(x), \quad 0 < x < \tau. \quad (8)$$

Proof. Choose $0 < \alpha < 1$ and M such that $\|f(y) - f(x)\| < M|y-x|^\alpha$.

Then

$$\int_0^\tau \|H(x-y)\{f(y) - f(x)\}\| dy < M \int_{-x}^{\tau-x} |z|^{\alpha-1} \|H(z)\| dz < \infty. \quad 12$$

As Ref. 12 shows that $N \stackrel{\text{def.}}{=} \int_{-\infty}^{\infty} |z|^{\alpha-1} \|H(z)\| dz < \infty$, the boundedness of F and its vanishing at ∞ are clear. Now define $G(x, y)$ such that

$$f(x) - f(y) = |x-y|^\alpha G(x, y).$$

Then G is bounded on $[0, \tau]^2$ (resp. $[0, \infty)^2$ for $\tau = \infty$) and, if necessary, by lowering α , one could choose G to be continuous.

Fix $\epsilon > 0$. If τ is finite, there exists a two-variable H -valued polynomial $p(x, y)$ such that the difference $\|G(x, y) - p(x, y)\| < (\epsilon/4N)$ ($0 \leq x, y \leq \tau$). At the same time the coefficients of the polynomial $p(x, y)$ could be taken to be polynomials in $H = L_2[-1, 1]$. Without any trouble one establishes that $q(x) = \int_0^\tau |x-y|^\alpha H(x-y) p(x, y) dy$ is continuous in x on $[0, \infty)$ with $\|q(x)\| \rightarrow 0$ as $x \rightarrow \infty$. Put $r(x) = \int_0^\tau H(x-y)\{f(y) - f(x)\} dy$. For $0 \leq x_1, x_2 < \infty$ we have the estimate

$$\|r(x_1) - r(x_2)\| < \|q(x_1) - q(x_2)\| + \epsilon \leq d|x_1 - x_2| + \frac{1}{2} \epsilon,$$

which, for finite τ , proves the continuity of r and thus of F . For infinite τ one proceeds differently. Singling out a compact interval $[a, b] \subseteq [0, \infty)$, one chooses finite σ such that

$$\left\| \int_\sigma^\infty H(x-y)\{f(y) - f(x)\} dy \right\| < \frac{1}{4} \epsilon \text{ for } a \leq x \leq b. \text{ Then one proves}$$

the continuity of $s(x) = \int_0^\sigma H(x-y)\{f(y) - f(x)\}dy$ on $[0, \infty)$, analogously to the continuity proof for finite τ , and establishes the continuity of F on $[0, \infty)$.

Let us derive (8). Premultiplying (8) by Q_\pm one gets the Cauchy problems

$$(TF_+)'(x) = -F_+(x) + Q_+f(x) \quad (0 < x < \tau), \quad (TF_+)(0) = 0; \quad (9a)$$

$$(TF_-)'(x) = -F_-(x) + Q_-f(x) \quad (0 < x < \tau), \quad (TF_-)(\tau) = 0. \quad (9b)$$

where $F_\pm(x) = Q_\pm F(x)$. Let us construct the solutions. Because in $\text{Ran} Q_\pm$ the self-adjoint operator $\mp T^{-1}Q_\pm$ generates a bounded analytic semi-¹³group and $Q_\pm f: (0, \tau) \rightarrow H_\pm$ are bounded and uniformly Hölder continuous, the above Cauchy problems both have unique solutions TF_+ and TF_- ¹⁴ and these solutions have the form

$$TF_+(x) = \int_0^x e^{-(x-y)T^{-1}} Q_+ f(y) dy, \quad TF_-(x) = -\int_x^\tau e^{-(x-y)T^{-1}} Q_- f(y) dy.$$

Define $G(x) = TF_+(x) + TF_-(x)$; then one can write

$$G(x) = T[e^{-xT^{-1}} Q_+ - e^{(\tau-x)T^{-1}} Q_-]f(x) + \int_0^\tau TH(x-y)\{f(y) - f(x)\}dy,$$

and thus $G(x) = TF(x)$ with $F(x)$ as in the statement of the lemma. Hence, TF is strongly differentiable on $(0, \tau)$ and satisfies Eq. (8).

Theorem 3. For finite (resp. infinite) τ let f be a bounded uniformly Hölder continuous function from $[0, \tau]$ (resp. $[0, \infty)$) into H . Then an essentially bounded (strongly measurable¹⁵) function $\psi: (0, \tau) \rightarrow H$ is a solution of the boundary value problem (5a)-(5b) (for finite τ) or (5a)-(5c) (for $\tau = \infty$), if and only if it is a solution of the vector-valued integral equation

$$\psi(x) - \xi \int_0^\tau c(y)H(x-y)B\psi(y)dy = \omega(x) \quad (0 < x < \tau), \quad (10)$$

where the right-hand side is given by

$$\begin{aligned} \omega(x) = & [e^{-xT^{-1}} Q_+ + e^{(\tau-x)T^{-1}} Q_-](\phi - f(x)) + f(x) + \\ & + \int_0^\tau H(x-y)\{f(y) - f(x)\}dy. \end{aligned} \quad (11)$$

Every such solution is continuous on $[0, \tau]$ for finite τ (resp. $[0, \infty)$ for $\tau = \infty$).

Proof. Note that

$$\omega(x) = [e^{-xT} Q_+^{-1} + e^{(\tau-x)T} Q_-^{-1}] \phi + F(x), \quad 0 < x < \tau,$$

which implies the differential equation

$$(T\omega)'(x) = -\omega(x) + f(x), \quad 0 < x < \tau.$$

The theorem can now be proved as is done for homogeneous media (i.e., $c(x) \equiv 1$).¹⁶ The boundedness of $\xi c(y)$ on $(0, \tau)$ assures the validity of all estimates made in the proof.

Using $e(\mu) \equiv 1$ and putting $\zeta(x) = \langle \psi(x), e \rangle$ one immediately derives from (10) and (7) the scalar integral equation

$$\zeta(x) - \frac{1}{2} \xi \int_0^\tau c(y) \text{Ei}(x-y) \zeta(y) dy = \langle \omega(x), e \rangle \quad (0 < x < \tau). \quad (12)$$

Conversely, once Eq. (12) is solved for $\zeta \in L_\infty(0, \tau)$, one puts

$$\psi(x) = \omega(x) + \frac{1}{2} \xi \int_0^\tau c(y) \zeta(y) H(x-y) e dy \quad (0 < x < \tau) \quad (13)$$

and gets an essentially bounded (strongly measurable) solution of Eq. (10). Hence, the boundary value problems (5a) - (5b) and (5a) - (5c) have been reduced to the scalar integral equation (12). For homogeneous media ($c(x) \equiv 1$) Eq. (12) was discovered in radiative transfer applications;¹⁷ inhomogeneous media versions appeared much later.¹⁸

Let us study Eq. (12) on the Banach space $BC(0, \tau)$. On this space the operator

$$(R_\zeta)(x) = \frac{1}{2} \int_0^\tau c(y) \text{Ei}(x-y) \zeta(y) dy \quad (14)$$

is bounded with norm

$$\|R\| \leq \frac{1}{2} \text{ess sup}_{0 < y < \tau} |c(y)| \cdot \int_{-\tau}^\tau \text{Ei}(z) dz, \quad (15)$$

and Eq. (12) is written in the form

$$(I-\xi R)\zeta = \bar{\omega}, \tag{16}$$

where $\bar{\omega}(x) = \langle \omega(x), e \rangle (0 < x < \tau)$. We shall come back to Eq. (16) in Section 3.

2. Operators leaving invariant a cone in a Banach space^{9,10}

Let X be a real Banach space. A non-empty closed subset K in X is called a cone in X if it has the following properties:

- (1) $x, y \in K \Rightarrow x+y \in K$, (2) $x \in K, \lambda > 0 \Rightarrow \lambda x \in K$, (3) $x, -x \in K \Rightarrow x = 0$.

A cone K in X is called reproducing if every $z \in X$ can be written in the form $z = x-y$ with $x, y \in K$. A cone K in X is called normal if there exists a constant M such that $x \in K$ and $y-x \in K$ imply that $0 \leq \|x\| \leq M \|y\|$.¹⁹ The set of nonnegative functions in $L_p[-1,1]$ ($1 \leq p < \infty$) is a normal and reproducing cone. The set of nonnegative functions in $BC(0, \tau)$ is a normal and reproducing cone also.

A bounded linear operator R on X is called positive (with respect to K) if R leaves invariant the cone K . For $0 \neq e \in K$ one defines R to be e-bounded below (resp. above) if for all $0 \neq x \in K$ there are $m \in \mathbb{N}$ and $\alpha > 0$ (resp. $m \in \mathbb{N}$ and $\beta > 0$) such that $R^m x - \alpha e \in K$ (resp. $\beta e - R^m x \in K$). If R is both e-bounded above and e-bounded below for the same $0 \neq e \in K$, for every $0 \neq x \in K$ there exist $m \in \mathbb{N}$ and $\alpha, \beta > 0$ such that $\{R^m x - \alpha e, \beta e - R^m x\} \subseteq K$, and in this case R is said to be e-positive.

When referring to spectral properties of R on real Banach spaces X , one complexifies X by embedding X into the complex space $\tilde{X} = X+iX$ and defining \tilde{R} by $\tilde{R}(x+iy) = Rx + iRy$. The complex space \tilde{X} will be a Banach space with respect to the norm

$$\|x+iy\| = \sup_{0 \leq \theta < 2\pi} \{\|x\| \cos \theta + \|y\| \sin \theta\} \quad (x, y \in X),$$

and in this way X is isometrically embedded in \tilde{X} and $\|\tilde{R}\| = \|R\|$. The spectral properties of R , such as eigenvalues and spectral radius, are simply defined as the ones of \tilde{R} . In particular, for the spectral radius one has $r(R) \stackrel{\text{def}}{=} r(\tilde{R}) = \lim_{n \rightarrow \infty} \|R^n\|^{1/n}$.

Proposition 1.²⁰ Let K be a normal and reproducing cone in X and R a positive operator on X . Then the spectral radius $r(R)$ of R belongs to the spectrum of R . In particular, R does not have eigenvalues λ such that $|\lambda| > r(R)$.

Proposition 2.^{9,10} Let K be a reproducing cone in X , and let R be a compact positive operator. Then either $r(R) = \{0\}$ or there exist $\lambda > 0$ and $0 \neq x \in K$ with $Rx = \lambda x$.

Proposition 3.^{9,10} Let K be a reproducing cone in X , and let R be e -positive for some $0 \neq e \in K$. Assume that $\tilde{R}x = \lambda x$ for some $\lambda \in \mathbb{C}$ and $0 \neq x \in K$. Then

- (1) the operator R is x -positive;
- (2) λ is a positive eigenvalue of R of geometric multiplicity one;
- (3) any other eigenvalue of \tilde{R} does not correspond to eigenvectors in K ;
- (4) all eigenvalues of \tilde{R} are contained in $\{\lambda\} \cup \{z \in \mathbb{C} : 0 < |z| < \lambda\}$.

Proposition 4.¹⁰ Let K be a reproducing cone in X , and let R be e -positive for some $0 \neq e \in K$. If λ is a positive eigenvalue of R corresponding to an eigenvector in K , then the equation

$$(\mu - R)x = y \quad (17)$$

does not have any solutions x in K whenever $0 \neq y \in K$ and $\mu \leq \lambda$.

Under the conditions of Proposition 4, certainly $\lambda \leq r(R)$. For $\mu > r(R)$ the equation (17) has a unique solution x for $y \in X$, and

$$x = (\mu - R)^{-1}y = \sum_{n=0}^{\infty} \mu^{-(n+1)} R^n y.$$

Thus $x \in K$ whenever $y \in K$.

3. Proofs of Theorems 1 and 2.

In Section 1 we have written Eq. (12) in the form (16), where R is given by (14). On the real Banach space $BC(0, \tau)$ we consider the cone K_c of nonnegative functions. Directly from the definitions one shows K_c to be a reproducing and normal cone invariant under R . Since one could factorize $R = CM$, where $(C\zeta)(x) = \frac{1}{2} \int_0^\tau E_i(x-y)\zeta(y)dy$ is a convolution operator²¹ and $(M\zeta)(x) = c(x)\zeta(x)$ is a multiplication by an L_∞ -function, the operator $L = CM$ is bounded on $L_\infty(0, \tau)$ with range contained in $BC(0, \tau)$; for finite τ , C is a compact operator from $L_\infty(0, \tau)$ into $BC(0, \tau)$, and thus $L = CM$ is compact on $BC(0, \tau)$.

LEMMA 2. Let $c(x) > 0$, except on a subset of $(0, \tau)$ of measure

zero. Then, for $e(\mu) \equiv 1$ in K_c and finite τ , the operator R is e -positive.

Proof. Certainly, if $\zeta \in K_c$, then for $0 < x < \tau$

$$0 \leq (R\zeta)(x) \leq \sup\{|\mathcal{R}\zeta(z)| : 0 < z < \tau\} \leq \|R\| \|\zeta\|,$$

and thus $\|R\| \cdot \|\zeta\| e - \zeta \in K_c$, thereby proving that R is e -bounded above.

Take $0 \neq \zeta \in K_c$, and choose $\varepsilon > 0$ and $(a,b) \subseteq (0,\tau)$ such that $\zeta(x) > \varepsilon$ on (a,b) . Because $\{x \in (0,\tau) : c(x) = 0\}$ has measure zero, there exist $E \subseteq (a,b)$ with positive measure and $\delta > 0$ such that $c(x) \geq \delta$ almost everywhere on E . So $c(x)\zeta(x) \geq \delta \varepsilon > 0$ almost everywhere on E . Thus

$$(R\zeta)(x) \geq \frac{1}{2} \delta \varepsilon \int_E E_i(x-y) dy = \frac{1}{2} \delta \varepsilon \int_{E-x} E_i(z) dz \geq M > 0,$$

which for finite τ implies that R is e -bounded below.

If there were an interval (a,b) in $(0,\tau)$ on which $c(x) = 0$, then $R\zeta = 0$ for all ζ which have their support on $[a,b]$. So in this case R cannot be e -positive.

Proof of Theorem 1. Define N by $(N\zeta)(x) = \sqrt{c(x)} \zeta(x)$. Then $N^2 = M$ and $R = CN^2$. The compact operators R on $BC(0,\tau)$, R on $L_2(0,\tau)$ and NCN on $L_2(0,\tau)$ have the same non-zero eigenvalues.²² As NCN is self-adjoint and non-zero, the spectral radius of NCN must be positive. Thus $r(R) > 0$. By Proposition 2, there exist $\lambda > 0$ and $0 \neq \eta \in K_c$ such that $R\eta = \lambda\eta$.

Because R is e -positive, Proposition 3 implies that λ is a simple eigenvalue of R and that R does not have any other eigenvectors in K_c than positive scalar multiples of η . Further, $\lambda = r(R) = r(NCN) = \|NCN\|$. Given τ , put $\xi_0(\tau) = \lambda^{-1}$.

First let $0 < \xi < \xi_0(\tau)$ (τ fixed). Then $|\xi| < r(R)^{-1}$, and thus Eq. (16) has a unique solution ζ for every $\tilde{\omega}$. Proposition 4 implies that $\zeta \in K_c$ whenever $\tilde{\omega} \in K_c$. However, if in the boundary value problem (5a)-(5b) the incoming flux $\phi \in K$ and the internal source term $f(x) \in K(0 < x < \tau)$, then (11) implies that the right-hand side $\omega(x)$ of the equivalent integral equation (10) assumes its values in K . By virtue of (12), $\tilde{\omega}(\cdot) = \langle \omega(\cdot), e \rangle = \int_{-1}^{+1} \omega(\cdot, \mu) d\mu \in K_c$,

and therefore $\zeta \in K_c$. Now (13) implies that ψ assumes its values in K . So we have established Part (a) of the theorem.

Let $\xi > \xi_0(\tau)$, and take $\phi \in K$ and $f(x) \in K(0 < x < \tau)$. Then, by (11), ω assumes its values in K . If there would be a solution $\psi: (0, \tau) \rightarrow K$ of (5a)-(5b), then $\zeta(\cdot) = \langle \psi(\cdot), e \rangle = \int_{-1}^{+1} \psi(\cdot, \mu) d\mu \in K_c$. Now Proposition 4 implies that $\zeta \equiv 0$, $\tilde{\omega} = \langle \omega(\cdot), e^{-1} \rangle \equiv 0$. Eq. (13) implies $\psi = \omega$. Because of (11) and the fact that f and ω assume their values in K , one has

$$\int_0^1 e^{-x/\mu} \phi(\mu) d\mu = 0, \quad \int_{-1}^0 e^{(\tau-x)/\mu} \phi(\mu) d\mu = 0, \quad \int_{-1}^{+1} F(x, \mu) d\mu = 0;$$

from this one easily derives $\phi = 0$ and $f \equiv 0$. Then $\omega \equiv 0$, and thus $\psi \equiv 0$. This proves Part (c) of the theorem.

Let $\xi = \xi_0(\tau)$. Then η satisfies Eq. (12) with zero right-hand side. Therefore, according to (13),

$$\psi(x) = \frac{1}{2} \xi_0(\tau) \int_0^\tau c(y) \eta(y) H(x-y) dy \quad (0 < x < \tau)$$

is a non-trivial solution of Eq. (10) with right-hand side $\omega \equiv 0$, which assumes its values in K . However, Theorem 3 implies that $\psi: (0, \tau) \rightarrow K$ is a non-trivial solution of Eqs. (5a)-(5b) with $\phi = 0$ and $f \equiv 0$. Now suppose, there is another non-trivial solution $\tilde{\psi}$, of Eqs. (5a)-(5b) (with $\phi = 0$, $f \equiv 0$); then this solution $\tilde{\psi}$ also satisfies Eqs. (10) (with $\omega \equiv 0$), and therefore, after proper multiplication by a positive constant, $\eta(x) = \langle \tilde{\psi}(x), e \rangle = \int_{-1}^{+1} \tilde{\psi}(x, \mu) d\mu$ ($0 < x < \tau$). Thus (13) implies that $\tilde{\psi} = \psi$. Hence, $\psi: (0, \tau) \rightarrow K$ is unique up to multiplication by a positive constant. Further, for $\phi \in K$ and $f(x) \in K$ ($0 < x < \tau$) there cannot be a solution $\psi: (0, \tau) \rightarrow K$ of (5a)-(5b), unless $\phi = 0$, $f \equiv 0$ and $\psi \equiv 0$. The latter is proved in exactly the same way as for $\xi > \xi_0(\tau)$. Thus Part (b) has been proved.

It remains to prove the monotonicity and continuity of the map $\tau \mapsto \xi_0(\tau)$, and the existence of $\lim_{\tau \downarrow 0} \xi_0(\tau)$ and $\lim_{\tau \rightarrow \infty} \xi_0(\tau)$. Fix $0 < \tau < \infty$. Let $\psi: (0, \tau) \rightarrow K$ be a non-trivial solution of the vector-valued integral equation

$$\psi(x) - \xi_0(\tau) \int_0^\tau c(y) H(x-y) B\psi(y) dy = 0, \quad 0 < x < \tau. \quad (18)$$

If $\psi(x_0) = 0$ for some $0 < x_0 < \tau$, then $c(y)H(x_0 - y)B\psi(y) \in K$ implies that $c(y)H(x_0 - y)B\psi(y) = 0$ for $0 < y < \tau$. As $c(y) > 0$, except on a set of measure zero, one gets $\langle \psi(y), e \rangle = 0$ for $0 < y < \tau$ (see (4b) for the form of B). Then (13) implies that $\psi \equiv 0$, which is a contradiction. Thus $\psi(x) > 0, 0 < x < \tau$. For $0 < x < \tau$ one gets from (18):

$$\psi(x) - \xi_0(\tau) \int_0^\sigma c(y)H(x-y)B\psi(y)dy = \xi_0(\tau) \int_\sigma^\tau c(y)H(x-y)B\psi(y)dy, \tag{19}$$

where the right-hand side is non-zero for every $0 < x < \sigma$.²³ As $\psi(x) \in K$ for $0 < x < \sigma$, the critical and supercritical cases are excluded, and thus

$$\xi_0(\tau) < \xi_0(\sigma)$$

(non-critical case). This establishes the monotonicity of $\tau \mapsto \xi_0(\tau)$. If this function would not be continuous, there would be a jump discontinuity at τ_0 . Put $\xi_+ = \lim_{\tau \downarrow \tau_0} \xi_0(\tau)$, $\xi_- = \lim_{\tau \uparrow \tau_0} \xi_0(\tau)$. Either $\xi_- > \xi_0(\tau) > \xi_+$ or $\xi_- > \xi_0(\tau) > \xi_+$. In the former case take $\xi_- > \xi_0(\tau) > \xi_+$, in the latter case $\xi_- > \xi_0(\tau) > \xi_+$. In both cases an argument as applied to prove monotonicity would lead to a contradiction. Thus $\tau \mapsto \xi_0(\tau)$ is continuous.

Finally, let us prove that $\lim_{\tau \rightarrow 0} \xi_0(\tau) = \infty$ and $\lim_{\tau \rightarrow \infty} \xi_0(\tau) = \xi_0(\infty) > 0$.

By virtue of (15), the norm of R , and thus its spectral radius $r(R)$ too, vanishes for $\tau \downarrow 0$, which shows that $\xi_0(\tau) = r(R)^{-1} \rightarrow \infty$ as $\tau \downarrow 0$. The same estimate (15) and the integrability of Ei on $(-\infty, \infty)$ ²⁴ show that for all finite τ

$$r(R) \leq \|R\| \leq \frac{1}{2} \|c\|_\infty \int_{-\infty}^{\infty} Ei(z) dz = \|c\|_\infty < \infty.$$

By the monotonicity of $\xi_0(\tau)^{-1} = r(R)$ as a function of τ , the existence of the limit $\lim_{\tau \rightarrow \infty} \xi_0(\tau) \stackrel{\text{def}}{=} \xi_0(\infty)$ is clear, and $\xi_0(\infty) > \|c\|_\infty^{-1} > 0$.

Proof of Theorem 2. The non-conservative case $\xi < \xi_0(\infty)$ is immediate from $\xi_0(\infty) = r(R)^{-1}$. The supercritical case $\xi > \xi_0(\infty)$ can be dealt with as follows: Consider Eq. (10) with $\tau = \infty$, $\omega(x) \in K(0 < x < \infty)$ and $\psi(x) \in K(0 < x < \infty)$, and fix finite τ with $\xi > \xi_0(\tau) > \xi_0(\infty)$. Writing

$$\psi(x) - \xi \int_0^{\tau} c(y)H(x-y)B\psi(y)dy = \omega(x) \quad (0 < x < \tau),$$

Theorem 1 implies that $\omega \equiv 0$ and $\psi \equiv 0$ on $(0, \tau)$. As this is the case for every finite τ with $\xi > \xi_0(\tau) > \xi_0(\infty)$, one has $\omega \equiv 0$ and $\psi \equiv 0$ on $(0, \infty)$, which establishes the theorem in the supercritical case.

The only situation not covered by Theorem 2 is the "conservative" half-space problem (5a)-(5c), where $\xi = \xi_0(\infty)$. This problem can be reduced to the integral equation

$$\zeta(x) - \frac{1}{2} \xi_0(\infty) \int_0^{\infty} c(y)Ei(x-y)\zeta(y)dy = \tilde{\omega}(x), \quad 0 < x < \infty, \quad (20)$$

where $\tilde{\omega}, \zeta \in BC(0, \infty)$. In homogeneous media ($c(y) \equiv 1, \xi_0(\infty) = 1$) and for $\tilde{\omega} = 0$ the above equation is the Schwarzschild-Milne integral equation,^{25,17} which does not have non-trivial bounded solutions.

For inhomogeneous media Eq. (20) with $\tilde{\omega} = 0$ has a solution

$0 \neq \zeta \in BC(0, \infty)$ if and only if $\xi_0(\infty)$ is an eigenvalue of R .

Tezcans⁷ reduced a criticality problem containing an anisotropy parameter β to a transport equation of the form

$$\mu \frac{\partial \psi'}{\partial x'} + \psi'(x', \mu') = \frac{1}{2} c' \int_{-1}^{+1} \psi'(x', \mu') d\mu', \quad -b < x' < b,$$

with boundary conditions

$$\psi'(-b, \mu) = 0 \quad (0 \leq \mu \leq 1), \quad \psi'(b, \mu) = 0 \quad (-1 \leq \mu \leq 0);$$

transforming $x'' = x' + b$ and $\tau = 2b$ one gets the boundary value problem (5a)-(5b) with $\xi = c', c(x) \equiv 1, f \equiv 0, \phi \equiv 0$. Tezcans⁷ parameters are related to the original collision ratio c , anisotropy parameter β , position (in units of neutron mean free path) x by

$$c' = c(1-\beta)/(1-c\beta), \quad x' = qx, \quad q = \sqrt{1-c^2\beta^2}.$$

Tezcans⁷ finding, which he explained physically and backed up numerically, is that the critical thickness vanishes for $\beta \uparrow (1/c)$. In our notation,

$$\lim_{\beta \uparrow (1/c)} \tau_0(c') / \sqrt{1-c^2\beta^2} = 0, \quad (21)$$

where $\xi_0(\tau_0(c)) = c$ for fixed $c > 1$. For homogeneous media ($c(x) \equiv 1$) one has to prove that

$$\lim_{\tau \downarrow 0} \tau^{-2} r(R_\tau)^{26} = \lim_{\tau \downarrow 0} (\tau^2 \xi_0(\tau))^{-1} = \infty, \tag{22}$$

because this would imply $\lim_{y \rightarrow \infty} \sqrt{y} \tau_0(y) = 0$, and thus (21). As R_τ is a self-adjoint operator on $L_2(0, \tau)$ and the spectra of the convolution operator R_τ on $L_2(0, \tau)$ and $BC(0, \tau)$ coincide,²¹ one has $r(R_\tau) = \|R_\tau\|_{L_2(0, \tau)} \geq \tau^{-1/2} \|R_\tau e\|_2$, where $e(x) \equiv 1$ has L_2 -norm $\sqrt{\tau}$. One easily computes that $(R_\tau e)(x) = \frac{1}{2} \{ \dot{1} - Ei_2(x) \} + \frac{1}{2} \{ \dot{1} - Ei_2(\tau - x) \}$ ²⁷ $\geq \frac{1}{2} \{ \dot{1} - Ei_2(\tau) \}$, $0 < x < \tau$, and thus

$$\tau^{-2} r(R_\tau) \geq \tau^{-5/2} \|R_\tau e\|_2 \geq \frac{1 - Ei_2(\tau)}{2\tau^2} = \frac{1}{2} \int_{\tau}^{\infty} \frac{1 - e^{-\omega}}{\omega^2} d\omega.$$

However, since the right-hand side tends to infinity as $\tau \downarrow 0$, Eq. (22) is clear. Herewith we have established (21).

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- (12) It is easy to compute that $\|H(x)\| = \sup \{ \mu^{-1} e^{-|x|/\mu} : 0 < \mu < 1 \}$, which equals $1/(ex)$ for $|x| < 1$ and $e^{-|x|}$ for $|x| > 1$.
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- (23) Note that $c(y)H(x-y)B\psi(y) \in K$, $0 < y < \sigma$. So $\int_{\sigma}^{\tau} c(y)H(x-y)B\psi(y)dy=0$ would imply $B\psi(y) = \langle \psi(y), e \rangle = 0$ ($0 < y < \tau$), and thus, by (13), $\psi \equiv 0$, which is a contradiction.
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