

The Hua equations imply (96). The converse question and a more explicit description of the kernel  $\nu^{-1}(0)$  are of considerable interest. Compare Ch. I, Theorem 3.23.

The dual transform to (89) would be

$$\check{\varphi}(gP_\Gamma) = \int_{P_\Gamma/L_o} \varphi(gqK_o) dq_{L_o} \quad \varphi \in \mathcal{D}(G_o/K_o).$$

The left invariant measure on  $L_oA_oN_o$  is just  $d\ell_o dadn$  and up to a constant

$$\int_X \varphi(x) dx = \int_{A_oN_o} \varphi(an \cdot o) dadn.$$

Hence the dual transform just maps  $\varphi$  into the constant function

$$\check{\varphi}(gP_\Gamma) = \int_X \varphi(x) dx.$$

Since

$$(\widehat{\nu}(D)\varphi)^\vee = \nu(D)\check{\varphi}$$

we see that if  $D$  is in the kernel of  $\nu$  then  $\widehat{\nu}(D)\mathcal{D}(G_o/K_o)$  is in the kernel of the dual transform  $\varphi \rightarrow \check{\varphi}$ . Here one could inquire about analogs of and perhaps sharpening of Ch. I, Theorem 3.24. A result in this spirit is the following result which is included among the exercises in this chapter:

If  $X = G/K$  is a symmetric space of rank one then

$$L_X(\mathcal{D}(X)) = \left\{ f \in \mathcal{D}(X) : \int_X f(x) dx = 0, \quad \int_X f(x) P(x, b) dx \equiv 0 \right\},$$

$P$  being the Poisson kernel on  $X \times B$ .

## §5. THE WAVE EQUATION ON SYMMETRIC SPACES.

### 1. INTRODUCTION. HUYGENS' PRINCIPLE.

It is a familiar phenomenon from daily life that waves propagate quite differently in 2 and 3 dimensions. When a pebble falls in water at a certain point  $P$ , circular ripples around  $P$  are formed. A given point  $Q$  near  $P$  will be hit by an initial ripple and later by residual waves.

In three dimensions the situation is quite different. A flash of light at a point  $P$  has an effect on the surface of a sphere around  $P$  after a certain time interval but then no more. There are no residual waves as those present on the water surface. The same is the case with sound waves; one has pure propagation without residual waves; thus music can exist in  $\mathbf{R}^3$  but not in  $\mathbf{R}^2$ .

There is a mathematical explanation of these phenomena. The propagation of waves in  $\mathbf{R}^n$  is governed by the wave equation

$$(1) \quad \frac{\partial^2 u}{\partial t^2} = \frac{\partial^2 u}{\partial x_1^2} + \cdots + \frac{\partial^2 u}{\partial x_n^2}, \quad u(x, 0) = f_o(x), \quad u_t(x, 0) = f_1(x).$$

Here  $u_t = \partial u / \partial t$ . The problem of solving (1) in terms of the initial data  $f_o, f_1 \in \mathcal{E}(\mathbf{R}^n)$  is a special case of the so called Cauchy problem for hyperbolic equations. Equation (1) can be solved by various methods; here we restate the solution obtained in [GGA], Ch. II, Exercise F1 by use of the Ásgeirsson mean value theorem. The result has different forms depending on the parity of  $n$ .

**Case I**  $n$  odd. With  $(M^r f)(x)$  denoting as usual the mean value of a function  $f$  on the sphere  $S_r(x)$  we put

$$(2) \quad (I_r f)(x) = \left( \frac{\partial}{\partial(r^2)} \right)^{\frac{n-3}{2}} (r^{n-2} (M^r f)(x)).$$

Then the solution to (1) is given by

$$(3) \quad u(x, t) = c_n \left[ \frac{\partial}{\partial t} (I_t f_o(x)) + (I_t f_1)(x) \right],$$

where  $c_n$  is the constant

$$c_n = \frac{\frac{1}{2} \Omega_n}{[\frac{1}{2}(n-3)]! \Omega_{n-1}}.$$

**Case II**  $n$  even. Here we put

$$(4) \quad (J_r f)(x) = \left( \frac{\partial}{\partial(r^2)} \right)^{\frac{n-2}{2}} (r^{n-2} (M^r f)(x)).$$

Then the solution to (1) is given by

$$(5) \quad u(x, t) = C_n \left[ \frac{\partial}{\partial t} \int_o^t r(t^2 - r^2)^{-\frac{1}{2}} (J_r f_o)(x) dr + \int_o^t r(t^2 - r^2)^{-\frac{1}{2}} (J_r f_1)(x) \right]$$

where  $C_n$  is the constant

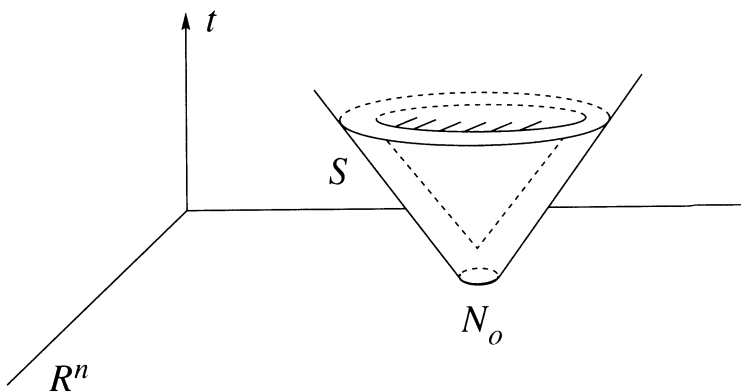
$$C_n = \frac{1}{[\frac{1}{2}(n-2)]!}.$$

In Case II formula (5) shows that  $u(x, t)$  is determined by the initial data  $f_o, f_1$  in a ball  $B_{t+\epsilon}(x)$  (for  $\epsilon > 0$  arbitrarily small). In Case I a stronger

statement holds:

- (6)  $u(x, t)$  is determined by the values of the initial data in an arbitrarily thin shell around  $S_t(x)$ .

To relate this better to light waves and sound waves suppose  $f_0$  and  $f_1$  have support in a small ball  $N_o = B_\epsilon(x_o)$  and consider the conical shell (see figure)



$S = \bigcup_{x \in N_o} C_x$  where  $C_x$  denotes the surface of the forward light cone with vertex  $x$ . Then Property (6) can also be stated in the equivalent form:

For all  $\epsilon > 0$  the support of the function  $(x, t) \rightarrow u(x, t)$  ( $t > 0$ ) is contained in the conical shell

$$(6') \quad S = \{(x, t) \mid t - \epsilon < |x_o - x| < t + \epsilon\},$$

whenever  $f_0$  and  $f_1$  have support in  $B_\epsilon(x_o)$ .

If (6') holds then for a fixed  $x \notin \overline{N_o}$  the function  $t \rightarrow u(x, t)$  is supported in a short time interval. This means pure propagation without residual effects.

For  $n = 2$  formula (5) only implies  $\text{supp}(u) \subset \tilde{S} = \bigcup_{x \in N_o} D_x$  ( $t > 0$ )

where  $D_x$  is the solid forward light cone with vertex  $x$ . Thus if  $x \notin \overline{N_o}$  is fixed, the function  $t \rightarrow u(x, t)$  will have support in an interval  $t \geq |x - x_o| - \epsilon$  as illustrated by the water waves.

Property (6) is an instance of the so called *Huygens' principle* (Hadamard's "Minor Premise"), which for general Riemannian manifolds  $X$  can be formulated as follows. Let  $L_X$  denote the Laplace operator on  $X$  and consider

the Cauchy problem,

$$(7) \quad \frac{\partial^2 u}{\partial t^2} = (L_X + c)u, \quad u(x, 0) = f_0(x), \quad u_t(x, 0) = f_1(x),$$

$c$  being some constant. *Huygens' principle* is said to hold for (7) if the solution  $u$  satisfies (6). Hadamard raised the question of finding all  $X$  for which (7) satisfies Huygens' principle. He proved that  $\dim X = \text{odd}$  is a necessary condition. More generally, the problem was raised with the operator  $\partial^2/\partial t^2 - L_X$  replaced with the Laplacian on a Lorentzian manifold and the surface  $t = 0$  by a spacelike hypersurface ([GGA], Ch. I, §6).

2. HUYGENS' PRINCIPLE FOR COMPACT GROUPS AND SYMMETRIC SPACES  $X = G/K$  ( $G$  COMPLEX).

To begin with let  $X = G/K$  be an arbitrary symmetric space of the noncompact type and  $\text{Exp}:\mathfrak{p} \rightarrow X$  the usual exponential mapping. As in [GGA], Ch. II, §3 let  $J$  denote the function on  $\mathfrak{p}$  defined by

$$(8) \quad \int_{G/K} f(x)dx = \int_{\mathfrak{p}} f(\text{Exp}X)J(X)dX$$

and put

$$(9) \quad \Omega(\text{Exp}Y) = J^{-\frac{1}{2}}(Y)L_{\mathfrak{p}}(J^{\frac{1}{2}}(Y)).$$

**Proposition 5.1.** *Let  $E$  be a  $K$ -invariant distribution on  $\mathfrak{p}$  and define  $\mathcal{E}_E \in \mathcal{D}'(X)$  by*

$$(10) \quad \mathcal{E}_E(f) = E(J^{\frac{1}{2}}(f \circ \text{Exp})), \quad f \in \mathcal{D}(X).$$

*Then*

$$(11) \quad (L_X + \Omega)\mathcal{E}_E = \mathcal{E}_{L_{\mathfrak{p}}E}.$$

*Proof.* We know from [GGA], Ch. II, §3 that

$$(12) \quad L_X^{\text{Exp}^{-1}}F = (L_{\mathfrak{p}} + \text{grad}(\log J))F$$

if  $F$  is a  $K$ -invariant function on  $\mathfrak{p}$ . Since

$$L_{\mathfrak{p}}(uv) = uL_{\mathfrak{p}}(v) + 2 \text{grad} u(v) + vL_{\mathfrak{p}}(u)$$

(12) can be written,

$$(L_X + \Omega)(F^{\text{Exp}}) = \left( J^{-\frac{1}{2}} L_{\mathfrak{p}}(J^{\frac{1}{2}} F) \right)^{\text{Exp}}.$$

Hence

$$\begin{aligned} ((L_X + \Omega)(\mathcal{E}_E))((J^{-\frac{1}{2}} F)^{\text{Exp}}) &= \mathcal{E}_E \left\{ (L_X + \Omega) \left( (J^{-\frac{1}{2}} F)^{\text{Exp}} \right) \right\} \\ &= \mathcal{E}_E \left\{ \left( J^{-\frac{1}{2}} L_{\mathfrak{p}}(F) \right)^{\text{Exp}} \right\} \\ &= E \left\{ J^{\frac{1}{2}} (J^{-\frac{1}{2}} L_{\mathfrak{p}}(F))^{\text{Exp}} \circ \text{Exp} \right\} = E(L_{\mathfrak{p}}(F)) = \\ (L_{\mathfrak{p}} E)(F) &= (J^{\frac{1}{2}} (L_{\mathfrak{p}} E))(J^{-\frac{1}{2}} F) = (J^{\frac{1}{2}} L_{\mathfrak{p}} E)^{\text{Exp}} (J^{-\frac{1}{2}} F)^{\text{Exp}}. \end{aligned}$$

Since both  $(L_X + \Omega)(\mathcal{E}_E)$  and  $(J^{\frac{1}{2}} L_{\mathfrak{p}} E)^{\text{Exp}}$  are  $K$ -invariant the calculation above implies that they coincide. But if  $f \in \mathcal{D}(X)$  then

$$\left( J^{\frac{1}{2}} L_{\mathfrak{p}} E \right)^{\text{Exp}}(f) = (L_{\mathfrak{p}} E) \left( J^{\frac{1}{2}} (f \circ \text{Exp}) \right) = \mathcal{E}_{L_{\mathfrak{p}} E}(f)$$

proving the proposition.

Now suppose  $G$  is complex. Then the  $K$ -radial part of  $L_{\mathfrak{p}}$  ([GGA], Ch. II, §3) is given by

$$(13) \quad \Delta(L_{\mathfrak{p}}) = \pi^{-1} L_{\mathfrak{a}} \circ \pi,$$

where  $\pi = \prod_{\alpha \in \Sigma^+} \alpha$ . Also

$$(14) \quad J(H) = \prod_{\alpha \in \Sigma^+} (\sinh \alpha(H))^2 / \pi(H)^2.$$

Since

$$(15) \quad \prod_{\alpha \in \Sigma^+} 2 \sinh \alpha(H) = \sum_{s \in W} \epsilon(s) e^{s\rho(H)}$$

we thus deduce from (9), (13)–(14) that

$$\Omega(\exp H) = \langle \rho, \rho \rangle \quad H \in \mathfrak{a}$$

so by  $K$ -invariance,  $\Omega \equiv \langle \rho, \rho \rangle$ . Because of Prop. 5.1 we thus take  $c = \langle \rho, \rho \rangle$  in (7) and take the shifted wave equation

$$(16) \quad \frac{\partial^2}{\partial t^2} = (L_X + |\rho|^2) u, \quad u(x, 0) = f_0(x), \quad u_t(x, 0) = f_1(x)$$

as the natural one on  $X$ . (The so-called Klein-Gordon equation has opposite sign on the constant term).

**Proposition 5.2.** Define the distribution  $E_t$  as follows:

$$(17) \quad E_t(\varphi) = \begin{cases} c_n(I_t\varphi)(0), & n \text{ odd} \\ C_n \int_0^t r(t^2 - r^2)^{-\frac{1}{2}} (J_r\varphi)(0) dr, & n \text{ even} \end{cases}$$

where  $\varphi \in \mathcal{D}(\mathfrak{p})$  and  $n = \dim X$ . Then the solution to

$$(18) \quad \frac{\partial^2 u}{\partial t^2} = L_{\mathfrak{p}}u, \quad u(x, 0) = f_o(x), \quad u_t(x, 0) = f_1(x)$$

is given by

$$(19) \quad u(x, t) = (f_o * E'_t)(x) + (f_1 * E_t)(x).$$

Here  $*$  denotes convolution on  $\mathfrak{p}$  and  $'$  denotes  $d/dt$ .

*Proof.* This is just a restatement of (3) and (5).

In accordance with Prop. 5.1 we put

$$(20) \quad \mathcal{E}_t(f) = E_t(J^{\frac{1}{2}}(f \circ \text{Exp})) \quad f \in \mathcal{D}(X).$$

**Theorem 5.3.** Let  $X = G/K$  be a symmetric space of the noncompact type,  $G$  complex. The solution to the shifted wave-equation (16) with  $f_o, f_1 \in \mathcal{E}(X)$  is given by

$$(21) \quad u(x, t) = (f_o \times \mathcal{E}'_t)(x) + (f_1 \times \mathcal{E}_t)(x),$$

where  $\times$  denotes the convolution on  $X$ .

If  $\dim X$  is odd then Huygens' principle holds for (16).

*Proof.* Apply  $L_{\mathfrak{p}}$  as well as  $\partial^2/\partial t^2$  to (19) and equate the results. Since  $f_o$  and  $f_1$  are arbitrary, we deduce that

$$(22) \quad L_{\mathfrak{p}}E_t = \frac{\partial^2 E_t}{\partial t^2},$$

Thus, by Prop. 5.1,

$$(23) \quad \frac{\partial^2 \mathcal{E}_t}{\partial t^2} = (L_X + |\rho|^2)\mathcal{E}_t.$$

Applying  $L_X$  to (21) we get by [GGA], Ch. II, Theorem 5.5,

$$(L_X)_x(u(x, t)) = (f_o \times L_X \mathcal{E}'_t)(x) + (f_1 \times L_X \mathcal{E}_t)(x)$$

so by (23)

$$(L_X + |\rho|^2)_x (u(x, t)) = f_0 \times \frac{\partial^2 \mathcal{E}'_t}{\partial t^2} + f_1 \times \frac{\partial^2 \mathcal{E}_t}{\partial t^2} = \frac{\partial^2 u}{\partial t^2}.$$

Secondly

$$(24) \quad \lim_{t \rightarrow 0} \mathcal{E}'_t(f_0) = \lim_{t \rightarrow 0} E'_t \left( J^{\frac{1}{2}}(f_0 \circ \text{Exp}) \right) = f_0(o)$$

and since  $E'_t$  is  $K$ -invariant and invariant under the symmetry  $Y \rightarrow -Y$ , the lift  $\tilde{\mathcal{E}}'_t$  of  $\mathcal{E}'_t$  to  $G$  is invariant under the inversion  $g \rightarrow g^{-1}$ . Hence

$$(25) \quad (f_0 \times \mathcal{E}'_t)(g \cdot o) = \int f_0(g \cdot x) d\mathcal{E}'_t(x),$$

which by (24) has limit  $f_0(g \cdot o)$  as  $t \rightarrow 0$ . Similarly

$$(26) \quad f_0 \times \mathcal{E}''_t \rightarrow 0 \text{ and } f_1 \times \mathcal{E}_t \rightarrow 0$$

as  $t \rightarrow 0$ . This proves that (21) does indeed satisfy the shifted wave equation.

Finally, if  $\dim X$  is odd,  $E_t$  and  $E'_t$  have support contained in  $S_t(0)$  and by (20),  $\mathcal{E}_t$  and  $\mathcal{E}'_t$  have support in  $S_t(o)$ . Since

$$(27) \quad u(g \cdot o, t) = \int_X f_0(g \cdot x) d\mathcal{E}'_t(x) + \int_X f_1(g \cdot x) d\mathcal{E}_t(x),$$

the value  $u(g \cdot o, t)$  is determined by the values of  $f_0$  and  $f_1$  in an arbitrary shell around  $g \cdot S_t(o) = S_t(g \cdot o)$ . This concludes the proof.

**Remark.** The constant  $\langle \rho, \rho \rangle$  in (16) is given by

$$(28) \quad \langle \rho, \rho \rangle = \frac{1}{12} \dim X.$$

To see this we extend  $\mathfrak{a}$  to a Cartan subalgebra  $\mathfrak{h}^{\mathbb{C}}$  of the complexification  $\mathfrak{g}^{\mathbb{C}}$ . Let  $\Delta(\mathfrak{g}^{\mathbb{C}}, \mathfrak{h}^{\mathbb{C}})$  be the corresponding set of nonzero roots. We order the roots in such a way that if  $\beta$  is a positive root then the restriction  $\bar{\beta} = \beta|_{\mathfrak{a}}$  belongs to  $\Sigma^+$  (cf. [DS], Ch. VI, Lemma 6.1). Let  $\Delta_+$  be the set of positive elements in  $\Delta(\mathfrak{g}^{\mathbb{C}}, \mathfrak{h}^{\mathbb{C}})$  and let  $\tilde{\rho} = \frac{1}{2} \sum_{\beta \in \Delta_+} \beta$ . Extending the

Cartan involution  $\theta$  of  $\mathfrak{g}$  to a complex automorphism of  $\mathfrak{g}^{\mathbb{C}}$  we know that  $\theta$  leaves  $\mathfrak{h}^{\mathbb{C}}$  invariant ([DS], Ch. VI, lemma 3.2) and by the compatibility of the orderings,  $\beta \rightarrow -\theta\beta$  is a permutation of  $\Delta_+$ . Hence  $\tilde{\rho} = -\theta\tilde{\rho}$  so  $\tilde{\rho}$  vanishes on  $\mathfrak{h}^{\mathbb{C}} \cap \mathfrak{k}$  whence  $\tilde{\rho} = \rho$ . The Killing forms  $B$  and  $B^{\mathbb{C}}$  of  $\mathfrak{g}^{\mathbb{C}}$  considered as a real (respectively complex) Lie algebra are related by  $B = 2\text{Re}B^{\mathbb{C}}$ . On the other hand, if  $G$  is simple we have by the ‘‘strange formula’’ of Freudenthal-deVries (cf. [GGA], p. 544) that

$$(29) \quad B^{\mathbb{C}}(\tilde{\rho}, \tilde{\rho}) = \frac{1}{24} \dim_{\mathbb{C}} G = \frac{1}{24} \dim X$$

so

$$B(\rho, \rho) = 2B^{\mathbf{C}}(\tilde{\rho}, \tilde{\rho}) = \frac{1}{12} \dim X$$

if  $G$  is simple. In the general semisimple case,  $X$  is a product and both sides of this equation are additive. This proves (28) in general.

Next we shall prove an analog of Theorem 5.3 for compact semisimple Lie groups  $K$ . Let  $R > 0$  be such that the mapping  $\exp : \mathfrak{k} \rightarrow K$  is a diffeomorphism of  $B_R(0)$  onto  $B_R(e)$ . The group  $K$  is furnished with an invariant Riemannian metric given by the negative of the Killing form  $Q$  of  $\mathfrak{k}$ .

Again let  $J$  be the volume element ratio given by

$$\int_K f(k) dk = \int_{B_R(0)} f(\exp Y) J(Y) dY$$

and if  $E$  is a distribution on  $B_R(0) \subset \mathfrak{k}$  invariant under  $Ad(K)$  we put (as in (10))

$$\mathcal{E}_E(f) = E(J^{\frac{1}{2}}(f \circ \exp)) \quad f \in \mathcal{D}(B_R(e)).$$

In analogy with (20) we put

$$(30) \quad \mathcal{E}_t(f) = E_t(J^{\frac{1}{2}}(f \circ \exp)) \quad \text{for} \quad f \in \mathcal{D}(B_R(e)),$$

where  $E_t$  is defined in (17) and  $|t| < R$ .

**Theorem 5.4.** *Let  $K$  be a simply connected compact semisimple Lie group with the bi-invariant metric given by the negative of the Killing form. Then the solution to the shifted wave equation on  $K$*

$$(31) \quad \frac{\partial^2 u}{\partial t^2} = (L_K - \frac{\dim K}{24})u, \quad u(k, 0) = f_0(k), \quad u_t(k, 0) = f_1(k)$$

where  $f_0, f_1 \in \mathcal{D}(B_R(e))$  has the solution

$$(32) \quad u(k, t) = (f_0 * \mathcal{E}'_t)(k) + (f_1 * \mathcal{E}_t)(k),$$

where  $*$  denotes convolution on  $K$ .

If  $\dim K$  is odd, then Huygens' principle holds for (31).

*Proof.* Let  $\mathfrak{t} \subset \mathfrak{k}$  be a maximal abelian subalgebra,  $\mathfrak{t}^{\mathbf{C}} \subset \mathfrak{k}^{\mathbf{C}}$  the corresponding complexifications,  $\Delta(\mathfrak{k}^{\mathbf{C}}, \mathfrak{t}^{\mathbf{C}})$  the corresponding set of roots and  $2\rho$  the sum of the corresponding set of positive roots relative to a specific ordering. Under the action of  $K$  on itself by conjugation  $T = \exp \mathfrak{t}$  is a transversal submanifold in the sense of [GGA], Ch. II, Lemma 3.3. As a



consequence of *loc. cit.*, Prop. 3.12, the corresponding radial part of the Laplacian  $L_K$  is given by

$$(33) \quad \Delta(L_K) = \delta^{-\frac{1}{2}} L_T \circ \delta^{\frac{1}{2}} + |\rho|^2,$$

where

$$(34) \quad \delta^{\frac{1}{2}}(\exp H) = \sum_{s \in W} (\det s) e^{s\rho(H)} \quad (H \in \mathfrak{t}),$$

$W$  being the Weyl group  $W(\mathfrak{k}, \mathfrak{t})$ . On the other hand, under the adjoint action of  $K$  of  $\mathfrak{k}$  with  $\mathfrak{t}$  as transversal submanifold the radial part  $\Delta(L_{\mathfrak{t}})$  given by Harish-Chandra's formula

$$(35) \quad \delta(L_{\mathfrak{t}}) = \pi^{-1} L_{\mathfrak{t}} \circ \pi,$$

where  $\pi$  is the product of the positive roots. On the positive Weyl chamber in  $\mathfrak{t}$   $J$  equals the ratio  $(\delta \circ \exp)/\pi^2$ . We can now prove the following result.

**Lemma 5.5.** *Let  $F \in C_c^\infty(B_R(0))$  be invariant under  $\text{Ad}(K)$ . Then*

$$(36) \quad L_K^{\exp^{-1}} F = \left( J^{-\frac{1}{2}} L_{\mathfrak{t}} \circ J^{\frac{1}{2}} \right) F + |\rho|^2 F.$$

Also, if  $E \in \mathcal{D}'(B_R(0))$  is  $\text{Ad}(\mathfrak{k})$ -invariant, and

$$(37) \quad \mathcal{E}_E(f) = E(J^{\frac{1}{2}}(f \circ \exp)) \quad f \in \mathcal{D}(B_R(e))$$

then

$$(38) \quad (L_K - |\rho|^2) \mathcal{E}_E = \mathcal{E}_{L_{\mathfrak{t}} E} \quad \text{on } B_R(e).$$

*Proof.* For (36) it suffices to verify that both sides agree on  $\mathfrak{t}$ . If  $H \in \mathfrak{t}$  we obtain from (33) and (35) if  $f = F^{\exp}$

$$\begin{aligned} \left( L_K^{\exp^{-1}} F \right) (H) &= (L_K f)(\exp H) \\ &= \delta^{-\frac{1}{2}}(\exp H) L_T(\delta^{\frac{1}{2}} f)(\exp H) + |\rho|^2 f(\exp H). \\ \left( J^{-\frac{1}{2}} L_{\mathfrak{t}}(J^{\frac{1}{2}} F) \right) (H) &= J^{-\frac{1}{2}}(H) \pi^{-1}(H) \left( L_{\mathfrak{t}}(\pi J^{\frac{1}{2}} F) \right) (H) \\ &= \delta^{-\frac{1}{2}}(\exp H) L_{\mathfrak{t}} \left( (\delta^{\frac{1}{2}} \circ \exp)(f \circ \exp) \right) (H), \end{aligned}$$

so (36) follows since  $L_{\mathfrak{t}}(g \circ \exp) = L_T(g) \circ \exp$  for  $g \in \mathcal{D}(\mathfrak{t} \cap B_R(0))$ .

For (38) let  $f \in \mathcal{D}(B_R(e))$  be invariant under  $K$ -conjugation. Then using (36),

$$\begin{aligned} ((L_K - |\rho|^2) \mathcal{E}_E)(f) &= E \left( J^{\frac{1}{2}} (L_K - |\rho|^2) f \circ \exp \right) \\ &= E \left( L_{\mathfrak{t}}(J^{\frac{1}{2}}(f \circ \exp)) \right) = (\mathcal{E}_{L_{\mathfrak{t}} E})(f). \end{aligned}$$

Next we can apply (29) to the simple components of the complexification of  $K$  and deduce

$$|\rho|^2 = \frac{1}{24} \dim K.$$

To finish the proof of Theorem 5.4, we note that by (24) and (38)

$$\begin{aligned} \frac{\partial^2 \mathcal{E}_t(f)}{\partial t^2} &= \frac{\partial^2 E_t(J^{\frac{1}{2}}(f \circ \exp))}{\partial t^2} = (L_{\mathfrak{t}} E_t)(J^{\frac{1}{2}}(f \circ \exp)) \\ &= \mathcal{E}_{L_{\mathfrak{t}} E_t}(f) = ((L_K - |\rho|^2) \mathcal{E}_t)(f). \end{aligned}$$

Thus  $u(x, t)$  as defined by (32) is a solution to the shifted wave equation. The initial conditions in (31) are verified in the same way as in Theorem 5.3. Again  $\mathcal{E}_t$  and  $\mathcal{E}'_t$  are invariant under the inversion  $k \rightarrow k^{-1}$  so

$$u(k, t) = \int_K f_0(kh) d\mathcal{E}'_t(h) + \int_K f_1(kh) d\mathcal{E}_t(h).$$

If  $\dim K$  is odd,  $\mathcal{E}'_t$  and  $\mathcal{E}_t$  are supported in  $S_t(e)$  ( $t < R$ ). Thus if  $k$  is such that  $kS_t(e) \subset B_R(e)$ ,  $u(k, t)$  is determined by the restrictions of  $f_0$  and  $f_1$  to an arbitrarily thin shell around  $S_t(e)$ . With these restrictions on  $t$  and  $k$ , the Huygens' principle (6) is verified and Theorem 5.4 proved.

### 3. HUYGENS' PRINCIPLE AND CARTAN SUBGROUPS.

In this subsection we shall study the solution of the shifted wave equation

$$(39) \quad \frac{\partial^2 u}{\partial t^2} = (L_X + |\rho|^2) u, \quad u(x, 0) = f_0(x), \quad u_t(x, 0) = f_1(x)$$

for arbitrary symmetric spaces  $X = G/K$  of the noncompact type, particularly as regards Huygens' principle. While Theorem 5.3 resulted from the special structure of  $L_X$  for  $G$  complex here we shall use the Fourier transform on  $X$ . We shall temporarily assume  $f_0, f_1 \in \mathcal{D}(X)$ .

Since the function

$$e_{\lambda, b}(x) = e^{(i\lambda + \rho)(A(x, b))} \quad x \in X$$

satisfies

$$L_X e_{\lambda, b} = -(\langle \lambda, \lambda \rangle + \langle \rho, \rho \rangle) e_{\lambda, b}$$

we see that if  $f \in \mathcal{D}(X)$ ,

$$(L_X f)^\sim(\lambda, b) = -(\langle \lambda, \lambda \rangle + \langle \rho, \rho \rangle) \tilde{f}(\lambda, b).$$

Thus taking formally the Fourier transform of (39) we would get, assuming the Fourier transform

$$(40) \quad \tilde{u}(\lambda, b, t) = \int_X u(x, t) e^{(-i\lambda + \rho)(A(x, b))} dx$$

defined for  $u$  as well as for the derivatives of  $u$ ,

$$(41) \quad \tilde{u}_{tt}(\lambda, b, t) + \langle \lambda, \lambda \rangle \tilde{u}(\lambda, b, t) = 0.$$

For  $\lambda \in \mathfrak{a}^*$  this would give

$$(42) \quad \tilde{u}(\lambda, b, t) = \tilde{f}_0(\lambda, b) \cos |\lambda|t + \tilde{f}_1(\lambda, b) \frac{\sin |\lambda|t}{|\lambda|}.$$

**Lemma 5.6.** *The function  $\lambda \longrightarrow |\lambda|^{-1} \sin(|\lambda|t)$  on  $\mathbf{R}^n$  is the Fourier transform*

$$(43) \quad \frac{\sin(|\lambda|t)}{|\lambda|} = \int_{\mathbf{R}^n} e^{-i\langle x, \lambda \rangle} dT_t(x),$$

where  $T_t \in \mathcal{E}'(\mathbf{R}^n)$  has support in the closed ball  $|x| \leq t$ . If  $n$  is odd and  $> 1$  then  $\text{supp}(T_t) \subset S_t(0)$ .

*Proof.* While this could be proved directly (see Remark below) the lemma is here essentially contained in Prop. 5.2. In fact, take  $f_0 \equiv 0$  in (18) and take Fourier transform in the  $x$  variable. Then

$$\tilde{u}(\lambda, t) = \tilde{f}_1(\lambda) \frac{\sin |\lambda|t}{|\lambda|}$$

so comparing with (19) we get  $E_t = T_t$  whence the lemma.

**Remark.** For a more direct proof of Lemma 5.6 consider the integral over  $S_t(0) \subset \mathbf{R}^n$  as the distribution

$$\tau_t(\varphi) = \int_{S_t(0)} \varphi(x) dw_t(x) = t^{n-1} \int_{S_1(0)} \varphi(tw) dw, \quad dw = dw_1.$$

Then ([GGA], Introduction, Lemma 3.6) we have with  $r = t|x|$ ,

$$\begin{aligned} \frac{1}{t} \int_{\mathbf{R}^n} e^{-i\langle x, u \rangle} d\tau_t(u) &= t^{n-2} \int_{\mathbf{S}^{n-1}} e^{-i\langle x, w \rangle t} dw = ct^{n-2} \frac{J_{(n-2)/2}(t|x|)}{(t|x|)^{(n-2)/2}} \\ &= c|x|^{2-n} r^{(n-2)/2} J_{(n-2)/2}(r) \end{aligned}$$

with  $c = (2\pi)^{n/2}$ . If  $n$  is odd, say  $n = 2m + 3$  then  $(n-2)/2 = m + \frac{1}{2}$  and we have the classical formula

$$\left( \frac{d}{z dz} \right)^m (z^{m+\frac{1}{2}} J_{m+\frac{1}{2}}(z)) = z^{\frac{1}{2}} J_{\frac{1}{2}}(z) = \left( \frac{2}{\pi} \right)^{1/2} \sin z$$

(which follows easily from the power series expansion for  $J_k(z)$ ). With  $r = t|x|$  we have

$$\frac{d}{t dt} = |x|^2 \frac{d}{r dr}.$$

Thus, putting  $\sigma_t = \left(\frac{d}{t dt}\right)^m \left(\frac{1}{t} \tau_t\right)$  we obtain

$$\int_{\mathbf{R}^n} e^{-i\langle x, u \rangle} d\sigma_t(u) = \text{const.} \frac{\sin t|x|}{|x|}.$$

This readily implies (43) for  $n$  odd and  $> 1$ , even with  $T_t$  supported in the sphere  $S_t(0)$ . To prove (43) for  $n$  replaced by  $n - 1$  we use the “method of descent” and simply restrict  $\lambda$  to the subspace  $\mathbf{R}^{n-1} \times 0$ . If  $T_t^*$  denotes the distribution

$$f \longrightarrow \int_{\mathbf{R}^n} f(x_1, \dots, x_{n-1}) dT_t(x), \quad f \in \mathcal{E}(\mathbf{R}^{n-1}),$$

then  $T_t^*(f) = 0$  for  $\text{supp}(f)$  disjoint from the ball  $|x| \leq t$  in  $\mathbf{R}^{n-1}$ . But then if  $\mu \in \mathbf{R}^{n-1}$ ,

$$(44) \quad \frac{\sin |\mu|t}{|\mu|} = \int_{\mathbf{R}^{n-1}} e^{-i(x, \mu)} dT_t^*(x),$$

proving (43) in the even dimensional case too.

Because of Lemma 5.6 we deduce from Euclidean Fourier transform theory that the function

$$\lambda \longrightarrow |\lambda|^{-1} \sin(|\lambda|t)$$

on  $\mathfrak{a}^*$  extends to a holomorphic function on  $\mathfrak{a}_{\mathbf{C}}^*$  which is of exponential type and slow growth in the sense of (40) Ch. III, §5. As a result of Cor. 5.9 in Ch. III there exists a distribution  $\mathcal{E}_t \in \mathcal{E}'(X)$  with support in  $\overline{B_t(o)}$  such that for  $\lambda \in \mathfrak{a}_{\mathbf{C}}^*, b \in B$ ,

$$(45) \quad \frac{\sin(|\lambda|t)}{|\lambda|} = \int_X e^{(-i\lambda + \rho)(A(x, b))} d\mathcal{E}_t(x).$$

While we arrived at the “propagator”  $\mathcal{E}_t$  by assuming  $f_o$  and  $f_1$  of compact support we can now solve (39) without this assumption.

**Theorem 5.7.** *With  $f_0, f_1 \in \mathcal{E}(X)$  the function*

$$(46) \quad u(x, t) = (f_0 \times \mathcal{E}'_t)(x) + (f_1 \times \mathcal{E}_t)(x)$$

*is a solution of the shifted wave equation (39).*

If  $\dim X$  is odd and if  $G$  has all its Cartan subgroups conjugate, then Huygens' principle holds.

*Proof.* We have

$$(L_X u)(x, t) = (f_0 \times L_X \mathcal{E}'_t)(x) + (f_1 \times L_X \mathcal{E}_t)(x)$$

and by (45),  $\mathcal{E}_0 = 0$ ,  $\mathcal{E}'_0 = \delta$ . For (39) it thus remains to prove

$$(L_X + |\rho|^2)\mathcal{E}_t = \frac{\partial^2}{\partial t^2}\mathcal{E}_t.$$

Both sides give the same result when applied to the function  $e_{\lambda, b}$  (Prop. 3.14, Ch. II) so they must coincide.

For the last statement of the theorem we first rewrite (43) for  $A$  :

$$(47) \quad \frac{\sin |\lambda|t}{|\lambda|} = \int_A e^{-i\lambda(\log a)} dT_t(a),$$

where  $\text{supp}(T_t) \subset \overline{B_t(e)}$  and  $\text{supp}(T_t) \subset S_t(e)$  if  $\dim A$  is odd and  $> 1$ .

**Lemma 5.8.** *In terms of the product decomposition  $\Xi = (K/M) \times A$  we have for the Radon transform  $\widehat{\mathcal{E}}_t$ ,*

$$\widehat{\mathcal{E}}_t = 1 \otimes e^\rho T_t.$$

*Proof.* Because of the uniqueness in (45),  $\mathcal{E}_t$  is invariant under  $K$ . Integrating (45) over  $b$  we get

$$(48) \quad \frac{\sin |\lambda|t}{|\lambda|} = \int_X \varphi_{-\lambda}(x) d\mathcal{E}_t(x),$$

where  $\varphi_{-\lambda}$  is the spherical function. If  $g = kan$  relative to the decomposition  $G = KAN$  we put as usual,  $a = \exp H(g)$ . Let  $\alpha \in \mathcal{E}(K/M)$ ,  $\beta(a) = e^{-(i\lambda+\rho)(\log a)}$  and  $\alpha^\natural$  the integral of  $\alpha$  over  $K/M$ . Then by the definition of  $\widehat{\mathcal{E}}_t$  (Ch. I, §3, (12)) and the  $K$ -invariance we have with  $(\alpha \otimes \beta)(kM, a) = \alpha(kM)\beta(a)$ ,

$$\begin{aligned} \widehat{\mathcal{E}}_t(\alpha \otimes \beta) &= \alpha^\natural \widehat{\mathcal{E}}_t(1 \otimes \beta) = \alpha^\natural \mathcal{E}_t((1 \otimes \beta)^\vee), \\ (1 \otimes \beta)^\vee(g \cdot o) &= \int_K (1 \otimes \beta)(gk \cdot \xi_o) dk = \int_K \beta(\exp H(gk)) dk = \varphi_{-\lambda}(g) \end{aligned}$$

so by (47)–(48),

$$\widehat{\mathcal{E}}_t(\alpha \otimes \beta) = \alpha^\natural \mathcal{E}_t(\varphi_{-\lambda}) = \alpha^\natural T_t(e^{-i\lambda}) = (1 \otimes e^\rho T_t)(\alpha \otimes \beta).$$

Since the functions  $\alpha \otimes \beta$  span a dense subspace of  $\mathcal{E}(\Xi)$  the lemma follows.

Let as in Ch. IV, §1,

$$\beta_R = \{\xi \in \Xi : d(o, \xi) < R\}, \quad \sigma_R = \{\xi \in \Xi : d(o, \xi) = R\}.$$

**Lemma 5.9.** *Suppose  $G$  has all its Cartan subgroups conjugate. Then if  $T \in \mathcal{E}'(X)$ ,*

$$\text{supp}(\widehat{T}) \subset \sigma_R \implies \text{supp}(T) \subset S_R(o).$$

*Proof.* From Ch. II, Theorem 3.13 we have the inversion formula for the Radon transform which we write in the form

$$(49) \quad f = (L\widehat{f})^\vee$$

and under our assumption,  $L$  is a differential operator on  $\Xi$ .

Let  $\epsilon > 0$  and suppose  $f$  satisfies  $\text{supp}(f) \subset B_{R-\epsilon}(o)$ . Then  $\text{supp}(\widehat{f}) \subset \bar{\beta}_{R-\epsilon}$  and since  $L$  in (49) is a differential operator,  $\text{supp}(L\widehat{f}) \subset \bar{\beta}_{R-\epsilon}$ . Hence

$$T(f) = T\left((L\widehat{f})^\vee\right) = \widehat{T}(L\widehat{f}) = 0$$

so  $\text{supp}(T) \cap B_{R-\epsilon}(o) = \emptyset$ . Since  $\epsilon$  is arbitrary,  $\text{supp}(T) \cap B_R(o) = \emptyset$ . On the other hand, by Ch. IV, Cor. 1.2,  $\text{supp}(T) \subset \overline{B_R(o)}$  so  $\text{supp}(T) \subset S_R(o)$  as desired.

The assumption on  $G$  implies that all restricted root subspaces are even-dimensional so since  $\dim X$  is odd,  $\dim A$  is odd. Assuming also  $\dim A > 1$ ,  $T_t$  in (47) has support in  $S_t(e)$  so by Lemma 5.8-5.9,  $\text{supp}(\mathcal{E}_t) \subset S_t(o)$ . Since  $\varphi_{-\lambda}(g \cdot o) = \varphi_\lambda(g^{-1} \cdot o)$  formula (48) implies that the lift of  $\mathcal{E}_t$  to  $G$  is invariant under  $g \rightarrow g^{-1}$ . Thus the solution formula (46) can be written

$$u(g \cdot o, t) = \int_X f_o(g \cdot x) d\mathcal{E}'_t(x) + \int_X f_1(g \cdot x) d\mathcal{E}_t(x)$$

and now Huygens' principle is obvious.

Finally, if  $X$  is of odd dimension and of rank one,  $X$  is an odd-dimensional hyperbolic space. In this case we have for the Euclidean Fourier transform

$$2 \frac{\sin |\lambda| t}{|\lambda|} = (\Phi_{B_{t(e)}})^\sim(\lambda), \quad 2 \cos |\lambda| t = (\Phi_{S_{t(e)}})^\sim(\lambda),$$

$\Phi_E$  denoting the characteristic function of a subset  $E \subset A$ , so the argument above gives Huygens' principle for (39) in the case when  $f_1 \equiv 0$ . However, it holds also in the general case because of the solution formula

$$(50) \quad u(x, t) = c_n \left[ \frac{\partial}{\partial t} (J_t f_o)(x) + (J_t f_1)(x) \right]$$

with  $c_n$  as in (3) and

$$(51) \quad (J_t f)(x) = \left( \frac{\partial}{\partial(2 \operatorname{ch} t)} \right)^{\frac{n-3}{2}} (\operatorname{sh}^{n-2} t (M^t f)(x)).$$

Here  $\mathbf{H}^n$  is taken with metric  $g$  with sectional curvature  $-1$  (cf. [GGA], pp. 343, 577), and the shifted wave equation considered is

$$(52) \quad \frac{\partial^2 v}{\partial s^2} = (L + (\frac{n-1}{2})^2) v, \quad v(x, 0) = v_o(x), \quad v_s(x, 0) = v_1(x).$$

For completeness this should be related to (39). We know from [DS], p. 566 that  $g = B/2(n-1)$ , where  $B$  is the Killing form. From [GGA], p. 582, we have for the single restricted root  $\alpha$ ,  $|\alpha|^2 = \frac{1}{2} m_\alpha^{-1} = (2(n-1))^{-1}$  so

$$|\rho|^2 = |\frac{1}{2}(n-1)\alpha|^2 = \frac{1}{8}(n-1).$$

Thus

$$L_X + |\rho|^2 = \frac{1}{2(n-1)}L + \frac{1}{8}(n-1) = \frac{1}{2(n-1)}(L + (\frac{n-1}{2})^2).$$

Thus the solutions to (39) and (52) correspond with

$$v(x, s) = u(x, (2(n-1))^{\frac{1}{2}}s).$$

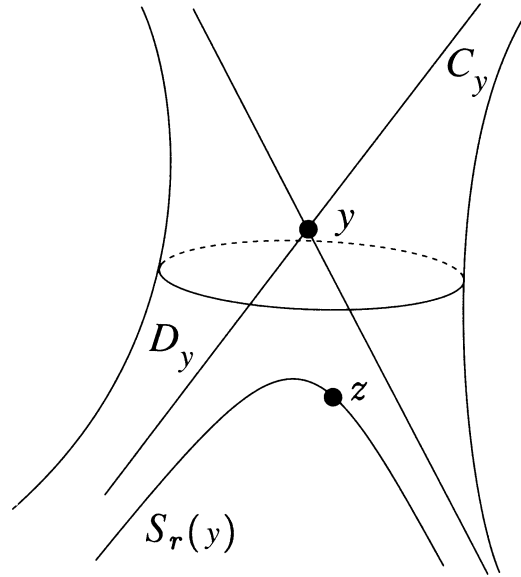
#### 4. ORBITAL INTEGRALS AND HUYGENS' PRINCIPLE.

Let  $X$  be an  $n$ -dimensional Lorentzian manifold, that is a pseudo-Riemannian manifold whose pseudo-Riemannian structure has signature  $(1, n-1)$ . We assume  $X$  has constant sectional curvature  $\varkappa = 0, -1$  or  $+1$  and therefore locally isometric to the respective spaces

$$(53) \quad \mathbf{O}_o(1, n-1)\mathbf{R}^n/\mathbf{O}_o(1, n-1), \mathbf{O}_o(1, n)/\mathbf{O}_o(1, n-1), \mathbf{O}_o(2, n-1)/\mathbf{O}_o(1, n-1),$$

the subscript  $o$  denoting identity component (cf. [GGA], Ch. I, §6).

Let  $X = G/H$  be one of the spaces (53). For  $y \in X$  let  $\mathbf{C}_y \subset X$  be the *light cone* with vertex  $y$  and  $\mathbf{D}_y \subset \mathbf{C}_y$  the *retrograde cone* with vertex  $y$ . Note that in [GGA], I, §6  $\mathbf{D}_y$  denotes the solid retrograde cone but here  $\mathbf{D}_y$  denotes the boundary of the solid retrograde cone. On the figure we consider the case  $\mathbf{O}_o(1, 2)/\mathbf{O}_o(1, 1)$  which is identified with the hyperboloid of one sheet in  $\mathbf{R}^3$ .



The group  $G = \mathbf{O}_o(1, 2)$  acts transitively on the hyperboloid. Here the light cone  $\mathbf{C}_y$  consists of two lines through  $y$  and  $\mathbf{D}_y$  of two half lines.

If  $n > 2$  the isotropy subgroup of  $G$  at  $y$  acts transitively on the connected set  $\mathbf{D}_y - \{y\}$ . Let  $o$  denote the origin  $eH$  in  $X, Y_1, \dots, Y_n$  a basis of the tangent space  $(G/H)_o$  in which the Lorentzian structure  $g_o$  has the form

$$(54) \quad g_o(Y, Y) = y_1^2 - \dots - y_n^2 \quad \text{if} \quad Y = \sum y_i Y_i.$$

Then if  $\text{Exp} : (G/H)_o \rightarrow G/H$  is the usual exponential mapping the retrograde cone  $\mathbf{D}_o$  is given by  $\mathbf{D}_o = \text{Exp}(D_o)$  where  $D_o$  is the retrograde cone in  $(G/H)_o$  :

$$y_1^2 - y_2^2 - \dots - y_n^2 = 0, \quad y_1 \leq 0.$$

The component of the hyperboloid  $g_o(Y, Y) = r^2$  which lies inside  $D_o$  is denoted  $S_r(o)$ ; if  $y \in X$  and if  $g \in G$  is chosen such that  $g \cdot o = y$  we put  $S_r(y) = g \cdot S_r(o)$ . This is a valid definition because  $h \cdot D_o \subset D_o$  for  $h \in H$  by the connectedness of  $H$ . Finally, we put

$$(55) \quad \mathbf{S}_r(y) = \text{Exp } S_r(y)$$



for  $r > 0$ ; for the last case in (53) we limit  $r$  to the interval  $(0, \pi)$  in order for  $\text{Exp}$  to be injective (see [GGA], I, §6).

Let  $dh$  denote the bi-invariant measure on the unimodular group  $H$ . Let  $y \in X, r > 0$  and select  $g \in G$  such that  $g \cdot o = y$  and let  $x \in \mathbf{S}_r(o)$ . We can then define the operator  $M^r$  (the orbital integral) by

$$(56) \quad (M^r u)(y) = \int_H u(gh \cdot x) dh, \quad u \in C_c(X).$$

As shown *loc. cit.*, Lemma 6.11, the measure  $dh$  can be normalized such that

$$(57) \quad (M^r u)(y) = \frac{1}{A(r)} \int_{\mathbf{S}_r(y)} u(z) d\mathbf{w}_r(z)$$

where  $d\mathbf{w}_r$  is the Riemannian measure on  $\mathbf{S}_r(y)$  and

$$(58) \quad A(r) = r^{n-1}, \quad (\sinh r)^{n-1}, \quad (\sin r)^{n-1}$$

in the three cases in (53). We then have the following result (*loc. cit.* §6).

**Theorem 5.10.** *Let  $X$  be one of the Lorentzian manifolds (53) of constant curvature  $\varkappa = 0, -1, +1$  and dimension  $n > 2$ .*

(i) *If  $v \in C_c(X)$  then*

$$(59) \quad a = \lim_{r \rightarrow 0} r^{n-2} (M^r v)(y) \text{ exists and is } \neq 0.$$

(ii)  *$M^r(Lu) = L_r(M^r u)$   $u \in \mathcal{D}(X)$  where  $L_r = d^2/dr^2 + A'(r)/A(r)d/dr$  is the radial part of the Laplacian  $L$ .*

(iii) *Suppose  $n$  is even,  $n = 2m$ , and put*

$$Q(L) = (L - \varkappa(n-3)2)(L - \varkappa(n-5)4) \cdots (L - \varkappa(1)(n-2)).$$

*Then if  $u \in \mathcal{D}(X)$ ,*

$$(60) \quad u(y) = c \lim_{r \rightarrow 0} r^{n-2} Q(L_r)((M^r u)(y)) \quad (n > 2)$$

*where the constant  $c$  is given by*

$$c^{-1} = (4\pi)^{m-1} (m-2)!.$$

Also if  $f$  is such that  $f \circ \text{Exp}$  has support in a ball  $B_A : \sum y_i^2 < A^2 < 1$  in  $X_o$  then (*loc. cit.* §6, No. 4) we have

$$(61) \quad |r^{n-2} (M^r f)(o)| \leq CA \sup |f|,$$

where  $C$  is a constant independent of  $r$ . If  $\varphi$  is the “smoothed-out” characteristic function of  $B_A$  we write with  $u \in \mathcal{D}(X)$ ,  $u_1 = u\varphi$ ,  $u_2 = u(1 - \varphi)$

$$(62) \quad r^{n-2}(M^r u)(o) = r^{n-2}(M^r u_1)(o) + r^{n-2}(M^r u_2)(o).$$

Take  $y = o$  in (59) and  $v \in C_c(X - \{o\})$ . The limit (59) is a positive, linear and  $H$ -invariant functional on  $C_c(X - \{o\})$  with support in the  $H$ -orbit  $\mathbf{D}_o - \{o\}$ ; hence it is “the”  $H$ -invariant measure on  $\mathbf{D}_o - \{o\}$  and we denote it by  $\mu$ . Then we have

$$\begin{aligned} & |r^{n-2}(M^r u)(o) - \int_{\mathbf{D}_o} u(z) d\mu(z)| \\ & \leq |r^{n-2}(M^r u_1)(o) - \int_{\mathbf{D}_o} u_1(z) d\mu(z)| + |r^{n-2}(M^r u_2)(o) - \int_{\mathbf{D}_o} u_2(z) d\mu(z)|. \end{aligned}$$

By (61) the first term on the right is  $O(A)$  uniformly in  $r$  and the second tends to 0 as  $r \rightarrow 0$ . Since  $A$  is arbitrary we conclude

$$(63) \quad \lim_{r \rightarrow o} r^{n-2}(M^r u)(o) = \int_{\mathbf{D}_o} u(z) d\mu(z).$$

**Corollary 5.11.** *Let  $n > 2$  be even and  $\delta$  the delta distribution at  $o$ .*

*Then*

$$(64) \quad \delta = cQ(L)\mu$$

In fact, using (ii) and (63) equation (60) becomes

$$u(o) = c \lim_{r \rightarrow o} r^{n-2}(M^r(Q(L)u))(o) = c \int_{\mathbf{D}_o} (Q(L)u)(z) d\mu(z),$$

which is (64).

Note that by (57), formula (63) can also be written

$$(63') \quad \lim_{r \rightarrow o} \frac{1}{r} \mathbf{w}_r = \mu.$$

**Remark.** Formula (64) shows that each factor

$$L_m = L - \varkappa(n - m)(m - 1) \quad (m = 3, 5, \dots, n - 1)$$

in  $Q(L)$  has a fundamental solution supported on the retrograde cone  $\mathbf{D}_o$ . It is known that this is equivalent to Huygens’ principle for the Cauchy

problem for the equation  $L_m u = 0$  with initial data given on some spacelike hypersurface (see Günther [1991] and [1988], Ch. IV, Cor. 1.13).

For the operator  $L - \frac{\kappa}{4}n(n - 2)$  which is  $L_m$  for  $m = 1 + 2 \lfloor \frac{n}{4} \rfloor$  this is *a priori* clear by the conformal invariance of Huygens' principle (Ørsted [1981]). Our space  $X$  is in fact conformally equivalent to  $\mathbf{R}^n$  so by the classical transformation formula for the Laplacian under conformal diffeomorphisms (see e.g. [GGA], p. 332) the null spaces of  $L_m$  and  $L_{\mathbf{R}^n}$  correspond. Huygens' principle for  $\mathcal{L}_{\mathbf{R}^n}$  explained in No. 1 of this section thus implies the same for  $L_m$  on  $X$ .

5. ENERGY EQUIPARTION.

Let  $X$  be an oriented manifold with Riemannian structure  $\langle, \rangle$  and let  $L_X$  denote the Laplace-Beltrami operator on  $X$ . For a vector field  $Y$  on  $X$  let  $\omega_Y$  be the 1-form given by  $\omega_Y(Z) = \langle Y, Z \rangle$  and extend  $\langle, \rangle$  to 1-forms by  $\langle \omega_Y, \omega_Z \rangle = \langle Y, Z \rangle$ . The star operator  $*$  on the Grassmann algebra  $\mathfrak{A}(X)$  maps  $p$ -forms to  $(n - p)$ -forms (if  $n = \dim X$ ), it operates pointwise, i.e.  $*(f\omega) = f(*\omega)$  if  $f \in \mathcal{E}(X)$  and it has the property that  $*(\omega_{i_1} \wedge \cdots \wedge \omega_{i_p}) = \pm \omega_{j_1} \wedge \cdots \wedge \omega_{j_{n-p}}$  if  $(\omega_i)$  is an orthonormal basis of 1-forms and  $\{i_1, \dots, i_p, j_1, \dots, j_{n-p}\}$  is a permutation of  $\{1, \dots, n\}$ , the sign being  $+$  or  $-$  depending on whether the permutation is even or odd (cf. [DS], p. 142 and [GGA], p. 330). The operator  $\delta$  is the linear operator on  $\mathfrak{A}(X)$  given by

$$\delta\omega = (-1)^{np+n+1} * d * \omega, \quad \omega \in \mathfrak{A}_p(X),$$

$d$  denoting as usual the exterior derivative. Then we know ([GGA], Ch. II, Ex. A.1) that for  $f \in \mathcal{E}(X)$ ,  $Y$  a vector field

$$(65) \quad \delta(\omega_Y) = -\operatorname{div} Y, \quad \omega_{\operatorname{grad} f} = df$$

and on  $\mathcal{E}(X)$ ,  $L_X = -\delta d$ . The familiar formula

$$\operatorname{div}(fY) = f \operatorname{div} Y + Yf$$

thus becomes

$$(66) \quad \delta(f\omega) = f\delta\omega - \langle df, \omega \rangle, \quad \omega \text{ 1-form, } f \in \mathcal{E}(X).$$

Consider now the shifted wave equation (7) with real valued initial data  $f_0$  and  $f_1$ . From general theory it is known that if  $f_0$  and  $f_1$  have compact support then so does the function  $x \rightarrow u(x, t)$ . More precisely, if  $\operatorname{supp}(f_i) \subset B_R(x_o)$  then for each  $t_o > 0$  the solution  $x \rightarrow u(x, t_o)$  has support in  $B_{R+t_o}(x_o)$ . Thus if  $f_0, f_1 \in \mathcal{D}(X)$  we can consider the *energy*

$$(67) \quad (\mathcal{E}_c u)(t) = \int_X (E_c u)(x, t) dx,$$

where

$$(68) \quad (E_c u)(x, t) = \frac{1}{2} (u_t^2 + \langle du, du \rangle - cu^2).$$

**Lemma 5.12.** (*Conservation of energy*).

$(\mathcal{E}_c u)(t)$  is independent of  $t$ ,

if  $u$  is a solution to (7).

*Proof.* From (68) we get

$$\begin{aligned} (E_c u)_t &= u_t u_{tt} + \langle du, du_t \rangle - cu u_t \\ &= u_t (u_{tt} - cu) + \langle du, du_t \rangle = u_t (-\delta du) + \langle du, du_t \rangle \\ &= -\delta(u_t du) \end{aligned}$$

by (66). Since by (65) this last expression is a divergence  $\operatorname{div} Z$  and since  $E_c u$  has compact support in the  $x$ -variable we have

$$\partial_t \int_X (E_c u)(x, t) dx = \int_X (E_c u)_t(x, t) dx = \int_X \operatorname{div} Z = 0,$$

proving the lemma.

We now specialize  $X$  to the spaces considered in Theorem 5.7 and take  $c = |\rho|^2$  so our shifted wave equation is

$$(69) \quad \frac{\partial^2 u}{\partial t^2} = (L_X + |\rho|^2)u, \quad u(x, 0) = f_0(x), \quad u_t(x, 0) = f_1(x)$$

with  $f_0, f_1, u$  real-valued and  $f_0, f_1 \in \mathcal{D}(X)$ . We split the energy  $\mathcal{E} = \mathcal{E}_{|\rho|^2}$  in two parts, the *kinetic energy* and the *potential energy*,

$$(70) \quad (\mathcal{K}u)(t) = \frac{1}{2} \int_X u_t^2 dx,$$

$$(71) \quad \mathcal{P}u(t) = \frac{1}{2} \int_X (\langle du, du \rangle - |\rho|^2 u^2) dx.$$

Since

$$\langle du, du \rangle = u \delta du - \delta(u du)$$

and since by (65),  $\int_X \delta(u du) dx = 0$  (71) can also be written

$$(72) \quad (\mathcal{P}u)(t) = -\frac{1}{2} \int_X u(L_X + |\rho|^2)u dx.$$

We can now state the energy equipartition theorem.

**Theorem 5.13.** *Let  $X$  be a symmetric space  $G/K$  of the noncompact type. We assume  $\dim X$  odd and that all Cartan subgroups of  $G$  are conjugate. Then if  $\text{supp } f_0, \text{supp } f_1 \subset B_R(o)$  we have*

$$\mathcal{K}u(t) = \mathcal{P}u(t) \quad \text{for } |t| \geq R.$$

*Proof.* Again we use the Fourier transform on  $X$  which, by (42), gives for the solution,

$$(73) \quad \tilde{u}(\lambda, b, t) = \Phi(\lambda, b) \cos pt + \Psi(\lambda, b) p^{-1} \sin pt,$$

where  $\Phi = \tilde{f}_0, \Psi = \tilde{f}_1, p = |\lambda|$ . Hence

$$\tilde{u}_t(\lambda, b, t) = -p \Phi(\lambda, b) \sin pt + \Psi(\lambda, b) \cos pt$$

so by the Plancherel formula,

$$\begin{aligned} 2w(\mathcal{K}u)(t) &= \int_{\mathfrak{a}^* \times B} \cos^2 pt |\Psi(\lambda, b)|^2 |\mathbf{c}(\lambda)|^{-2} d\lambda db \\ &\quad + \int_{\mathfrak{a}^* \times B} \sin^2 pt p^2 |\Phi(\lambda, b)|^2 |\mathbf{c}(\lambda)|^2 d\lambda db \\ &\quad - \int_{\mathfrak{a}^* \times B} \sin pt \cos pt p (\Phi \bar{\Psi} + \Psi \bar{\Phi})(\lambda, b) |\mathbf{c}(\lambda)|^{-2} d\lambda db. \end{aligned}$$

Hence, by the double-angle identities,

$$(74) \quad \begin{aligned} 2w(\mathcal{K}u)(t) &= \frac{1}{2} \|\Psi\|^2 + \frac{1}{2} \|p\Phi\|^2 \\ &\quad + \frac{1}{2} \int_{\mathfrak{a}^* \times B} \cos 2pt (|\Psi|^2 - |p\Phi|^2) |\mathbf{c}(\lambda)|^{-2} d\lambda db \\ &\quad - \frac{1}{2} \int_{\mathfrak{a}^* \times B} \sin 2pt p (\Phi \bar{\Psi} + \Psi \bar{\Phi}) |\mathbf{c}(\lambda)|^{-2} d\lambda db, \end{aligned}$$

the norms being taken in  $L^2(\mathfrak{a}^* \times B, |\mathbf{c}(\lambda)|^{-2} db d\lambda)$ . On the other hand, by (70), (72),

$$2(\mathcal{E}u)(t) = \int_X u_t^2 dx - \int_X u(L_X + |\rho|^2)u dx,$$

which by Lemma 5.12 and the Plancherel formula equals

$$2(\mathcal{E}u)(0) = (\|\Psi\|^2 + \|p\Phi\|^2)w^{-1}.$$

Since  $\mathcal{E}u(t) = (\mathcal{K}u)(t) + (\mathcal{P}u)(t)$  it will suffice to prove that

$$\mathcal{I}(t) = 0 \quad \text{for } |t| > R,$$

where  $\mathcal{I}(t)$  is the sum of the two last terms in (74). Actually each of the terms will be shown to vanish for  $|t| > R$ . Now we have by Ch. III, (25) in §5,

$$\begin{aligned} \Phi(\lambda, kM) &= \int_X f_0(x) e^{(\rho-i\lambda)(A(x, kM))} dx \\ &= \int_A \widehat{f}_0(ka \cdot \xi_o) e^{(\rho-i\lambda)(\log a)} da, \quad \lambda \in \mathfrak{a}^*. \end{aligned}$$

Here we put  $q = \lambda/p$  for the angular spherical coordinate in  $\mathfrak{a}^*$  and write

$$\begin{aligned} \Phi(p; q, kM) &= \Phi(\lambda, kM) = \int_X f_0(x) e^{(\rho-ipq)(A(x, kM))} dx \\ &= \int_A \widehat{f}_0(ka \cdot \xi_o) e^{(\rho-ipq)(\log a)} da. \end{aligned}$$

For  $(q, kM)$  fixed this defines an extension of  $p \rightarrow \Phi(p; q, kM)$  from  $\mathbf{R}^+$  to a holomorphic function  $\Phi^*(p, q, kM)$  of a single complex variable  $p = \sigma + i\tau$ . The last formula shows, since  $\widehat{f}_0(ka \cdot \xi_o) = 0$  for  $|\log a| > R$  and  $|q| \leq 1$ , that  $p \rightarrow \Phi(p; q, kM)$  is an entire function on  $\mathbf{C}$  of exponential type  $R$ , uniformly in  $(q, kM)$ . The complex conjugate  $\overline{\Phi}(\bullet, q, kM)$  has a similar extension to an entire function  $p \rightarrow \overline{\Phi}^*(p; q, kM)$  of exponential type  $R$ . Since  $f_0$  is real-valued we actually have

$$\overline{\Phi}^*(p; q, b) = \Phi^*(p; -q, b).$$

Similar construction is carried out for  $\Psi$ . Recall also ([GGA], Ch. IV, Cor. 6.15) that under our assumption on  $G$ ,  $\mathbf{c}(\lambda)^{-1}$  is a polynomial  $\gamma(\lambda)$  in  $i\lambda$  with real coefficients. We put  $\gamma(p; q) = \gamma(pq)$  for  $p \in \mathbf{R}^+, |q| = 1$ . Now

$$|\mathbf{c}(\lambda)|^{-2} = \mathbf{c}(\lambda)^{-1} \mathbf{c}(-\lambda)^{-1}, \quad \lambda \in \mathfrak{a}^*$$

so for a fixed  $q \in \mathfrak{a}^*, |q| = 1$ ,  $(|\mathbf{c}(\lambda)|^{-2})|_{\mathfrak{a}^*}$  has the holomorphic extension  $\gamma(p; q)\gamma(p; -q)$  to  $p \in \mathbf{C}$ . We now switch to polar coordinates  $(p, q)$  in the integrals in (74) so  $d\lambda = p^{\ell-1} dp dq$  where  $dq$  is the (suitably normalized) measure on the unit sphere  $S_1$  in  $\mathfrak{a}^*$ . Then

$$\begin{aligned} \mathcal{I}(t) &= \int_0^\infty (p^{\ell-1} \cos 2pt \int_{S_1 \times B} Y_o(p, q, b) db dq) dp \\ &\quad + \int_0^\infty (p^{\ell-1} \sin 2pt \int_{S_1 \times B} Y_1(p, q, b) db dq) dp, \end{aligned}$$

where

$$\begin{aligned} 2Y_0(p; q, b) &= \\ & [\Psi^*(p; q, b)\Psi^*(p, -q, b) - p^2\Phi^*(p; q, b)\Phi^*(p; -q, b)] \gamma(p; q)\gamma(p; -q) \\ 2Y_1(p; q, b) &= \\ & -p[\Phi^*(p; q, b)\Psi^*(p; -q, b) + \Phi^*(p; -q, b)\Psi^*(p; q, b)] \gamma(p; q)\gamma(p; -q). \end{aligned}$$

These functions  $Y_0, Y_1$  are for fixed  $(q, b)$  entire functions of  $p$  of exponential type  $2R$ . By the uniformity mentioned, the same is the case for the functions

$$Z_j(p) = p^{\ell-1} \int_{S_1 \times B} Y_j(p; q, b) db dq, \quad (j = 0, 1)$$

and the expression for  $\mathcal{I}(t)$  becomes

$$\mathcal{I}(t) = \int_0^\infty Z_0(p) \cos 2pt dp + \int_0^\infty Z_1(p) \sin 2pt dp.$$

On the other hand,  $\Phi^*(-p; q, b) = \Phi^*(p; -q, b)$  and similarly for  $\Psi^*$ . The assumptions on  $X$  and  $G$  imply that  $\dim A$  is odd so

$$Z_0(-p) = Z_0(p), \quad Z_1(-p) = -Z_1(p),$$

whence

$$\mathcal{I}(t) = \frac{1}{2} \int_{-\infty}^\infty (Z_0(p) + iZ_1(p)) e^{-2ipt} dp.$$

Here  $Z_0$  and  $Z_1$  are entire functions of exponential type  $2R$  so by the classical Paley-Wiener theorem  $\mathcal{I}(t) = 0$  for  $|2t| > 2R$ . This proves Theorem 5.13.

## 6. THE FLAT CASE REVISITED.

We shall now write down the solution to the Cauchy Problem (1) in a different form. We first discuss the Radon transform version and will use the inversion formula from Ch. I, §1, (9).

**Lemma 5.14.** *If  $h \in C^2(\mathbf{R})$  and  $\omega \in \mathbf{S}^{n-1}$  then the function*

$$(75) \quad v(x, t) = h(\langle x, \omega \rangle + t)$$

*satisfies  $Lv = \partial^2/\partial t^2 v$ .*

The proof is obvious. It is now easy on the basis of the inversion formula and (75) to write down the solution to the Cauchy problem (74).

**Theorem 5.15.** *The solution to (1) with  $f_0, f_1 \in \mathcal{D}(\mathbf{R}^n)$  is given by*

$$(76) \quad u(x, t) = \int_{\mathbf{S}^{n-1}} (Sf)(\omega, \langle x, \omega \rangle + t) d\omega,$$

where

$$(77) \quad Sf = \begin{cases} c(\partial^{n-1}\widehat{f}_0 + \partial^{n-2}\widehat{f}_1), & n \text{ odd, } > 1. \\ c[\mathcal{H}(\partial^{n-1}\widehat{f}_0 + \partial^{n-2}\widehat{f}_1)], & n \text{ even.} \end{cases}$$

Here  $\partial = \partial/\partial p$ , and the constant  $c$  equals

$$c = \frac{1}{2}(2\pi i)^{1-n}.$$

*Proof.* Because of Lemma 5.14 we just have to check the initial conditions in (1).

(i) If  $n > 1$  is odd then  $\omega \rightarrow (\partial^{n-1}\widehat{f}_0)(\omega, \langle x, \omega \rangle)$  is an even function on  $\mathbf{S}^{n-1}$  but  $\omega \rightarrow (\partial^{n-2}\widehat{f}_1)(\omega, \langle x, \omega \rangle)$  is odd. Thus  $u(x, 0) = f_0(x)$  by the inversion formula. Applying  $\partial/\partial t$  to the right hand side of (76) and putting  $t = 0$  gives  $u_t(x, 0) = f_1(x)$ , this time because the function  $\omega \rightarrow (\partial^n \widehat{f}_0)(\omega, \langle x, \omega \rangle)$  is odd and  $\omega \rightarrow (\partial^{n-1}\widehat{f}_1)(\omega, \langle x, \omega \rangle)$  is even.

(ii)  $n$  even. Here the proof is the same if one remarks that  $c\mathcal{H}$  interchanges even and odd functions on  $\mathbf{R}$ .

**Definition.** For the initial data  $f = \{f_0, f_1\}$  we shall refer to the function  $Sf$  in (77) on  $\mathbf{S}^{n-1} \times \mathbf{R}$  as the *source*.

**Definition.** (Lax–Phillips) The wave  $u(x, t)$  is

*outgoing* if  $u(x, t) = 0$  in the *forward cone*  $|x| < t$

*incoming* if  $u(x, t) = 0$  in the *backward cone*  $|x| < -t$ .

**Corollary 5.16.** *The solution  $u(x, t)$  to (1) is*

(i) *outgoing if and only if  $(Sf)(\omega, s) = 0$  for  $s > 0$ .*

(ii) *incoming if and only if  $(Sf)(\omega, s) = 0$  for  $s < 0$ .*

*Proof.* For (i) suppose  $(Sf)(\omega, s) = 0$  for  $s > 0$ . Then if  $|x| < t$  we have  $\langle x, \omega \rangle + t > -|x| + t > 0$  so by (76)  $u$  is outgoing. Conversely, suppose



$u(x, t) = 0$  for  $|x| < t$ . Let  $t_0 > 0$  be arbitrary and  $\varphi \in \mathcal{D}(t_0, \infty)$ . Then if  $|x| < t_0$  we have

$$\begin{aligned} 0 &= \int_{\mathbf{R}} u(x, t) \varphi(t) dt = \int_{\mathbf{S}^{n-1}} d\omega \int_{\mathbf{R}} (Sf)(\omega, \langle x, \omega \rangle + t) \varphi(t) dt \\ &= \int_{\mathbf{S}^{n-1}} d\omega \int_{\mathbf{R}} (Sf)(\omega, p) \varphi(p - \langle x, \omega \rangle) dp = 0. \end{aligned}$$

Taking  $\partial^{|k|} / \partial x_{i_1} \dots \partial x_{i_k}$  at  $x = 0$  we deduce

$$\int_{\mathbf{R}} \left( \int_{\mathbf{S}^{n-1}} (Sf)(\omega, p) \omega_{i_1} \dots \omega_{i_k} d\omega \right) (\partial^{|k|} \varphi)(p) dp = 0$$

for each  $k$  and each  $\varphi \in \mathcal{D}(t_0, \infty)$ . Integrating by parts in the  $\mathbf{R}$ -integral we deduce that the function

$$(78) \quad p \rightarrow \int_{\mathbf{S}^{n-1}} (Sf)(\omega, p) \omega_{i_1} \dots \omega_{i_k} d\omega$$

has its  $k$ th derivative  $\equiv 0$  for  $p > t_0$ . Having compact support if  $n$  is odd and being  $O(\log |p|)$  if  $n$  is even we deduce that the function (78) vanishes identically for  $p > t_0$ . Now varying  $k$  we see that  $(Sf)(\omega, p)$  vanishes for  $p > t_0$ . Since  $t_0 > 0$  was arbitrary this proves (i). Part (ii) is proved in the same way.

Since  $f_0$  and  $f_1$  have compact support we can rephrase the corollary as follows:

*For  $n$  odd the solution  $u(x, t)$  is*

- (i) *outgoing if and only if  $(\partial \widehat{f}_0 + \widehat{f}_1)(w, s) = 0$  for  $s > 0$ .*
- (ii) *incoming if and only if  $(\partial \widehat{f}_0 + \widehat{f}_1)(w, s) = 0$  for  $s < 0$ .*

The Cauchy Problem (1) can also be solved by using the Fourier transform

$$\widetilde{f}(\zeta) = \int_{\mathbf{R}^n} f(x) e^{-i(x, \zeta)} dx.$$

Assuming the function  $x \rightarrow u(x, t)$  in  $\mathcal{S}(\mathbf{R}^n)$  for a given  $t$  we obtain

$$\widetilde{u}_{tt}(\zeta, t) + (\zeta, \zeta) \widetilde{u}(\zeta, t) = 0.$$

Solving this ordinary differential equation with the initial data in (1) we get

$$\widetilde{u}(\zeta, t) = \widetilde{f}_0(\zeta) \cos(|\zeta|t) + \widetilde{f}_1(\zeta) \frac{\sin(|\zeta|t)}{|\zeta|}.$$

By Lemma 5.6 (and subsequent Remark) we deduce

$$(79) \quad u(x, t) = (f_0 * T'_t)(x) + (f_1 * T_t)(x), \quad T'_t = \frac{d}{dt}T_t$$

where  $T_t$  has support in  $B_{|t|}(0)$  and

$$(80) \quad \frac{\sin(|\zeta|t)}{|\zeta|} \int_{\mathbf{R}^n} e^{-\langle \zeta, x \rangle} dT_t(x), \quad \cos(|\zeta|t) = \int_{\mathbf{R}^n} e^{-\langle \zeta, x \rangle} dT'_t.$$

Formula (79) gives a solution to (1) for any  $f_0, f_1 \in \mathcal{E}(\mathbf{R}^n)$ .  
 Since

$$(81) \quad LT_t = \frac{d^2}{dt^2}T_t, \quad LT'_t = \frac{d^2}{dt^2}T'_t$$

$$(82) \quad (T_t)_{t=0} = 0, (T'_t)_{t=0} = \delta_0,$$

one refers to the pair  $\{T_t, T'_t\}$  as the *fundamental solution* to the wave equation.

For  $n$  odd, (76) implies immediately the classical Huygens' principle in the form (6').

In fact, for  $n$  odd (76) implies

$$(83) \quad u(0, t) = 0 \quad \text{for } |t| \geq R.$$

If  $y \in \mathbf{R}^n$  the translated initial data  $f_0^{-y}, f_1^{-y}$  give the solution  $x \rightarrow u(x + y, t)$ . Thus by (83) since  $\text{supp}(f_i^{-y}) \subset B_{R+|y|}(0)$ ,

$$u(y, t) = 0 \quad \text{for } t \geq R + |y|.$$

On the other hand, since  $T_t$  above has support in  $B_{|t|}(0)$ , (79) implies that  $u$  has support in the region  $|x| \leq |t| + R$ .

For  $n > 1$  odd one has also the following limit theorem of Friedlander.

**Theorem 5.17.** *For  $n$  odd the solution  $u$  to (1) satisfies*

$$\lim_{|t| \rightarrow \infty} t^{\frac{n-1}{2}} u_t((t+p)\omega, t) = (-\partial^{(n+1)/2} \widehat{f}_0 + \partial^{(n-1)/2} \widehat{f}_1)(\omega, p).$$

For a proof see the book Lax-Phillips, [15, p.108] which also contains the proof above for Cor. 5.16 for  $n$  odd.

### 7. THE MULTITEMPORAL WAVE EQUATION ON $X = G/K$

While the wave equation can be considered on any Riemannian manifold (see (7)), Semenov-Tjan-Shanski in [1976] introduced a kind of a

Cauchy problem on the symmetric space  $X = G/K$  in which the time variable is multidimensional.

For convenience we review some of the notation from Chapter III whose results will be extensively used. We have the space  $X = G/K$ , the Cartan decomposition  $\mathfrak{g} = \mathfrak{k} + \mathfrak{p}$ , the corresponding Cartan involution  $\theta$ , the Killing form  $B$  and we put  $\langle Y, Z \rangle = B(Y, Z)$  for  $Y, Z \in \mathfrak{p}$ . We fix the maximal abelian subspace  $\mathfrak{a} \subset \mathfrak{p}$ , its dual  $\mathfrak{a}^*$ ,  $\Sigma = \Sigma(\mathfrak{g}, \mathfrak{a})$  the system of roots of  $\mathfrak{g}$  relative to  $\mathfrak{a}$  and  $\Sigma_0$  the set of indivisible roots.

The Weyl group  $W$  of  $\Sigma$  acts on  $\mathfrak{a}$  and  $\mathfrak{a}^*$  by

$$(\sigma \cdot \lambda)(H) = \lambda(\sigma^{-1}H) = \lambda^\sigma(H) \quad \text{for } \sigma \in W, \lambda \in \mathfrak{a}^*, H \in \mathfrak{a}$$

and this action extends to automorphisms of the symmetric algebras  $S(\mathfrak{a})$  and  $S(\mathfrak{a}^*)$ .

For  $\lambda \in \mathfrak{a}$ , let  $H_\lambda \in \mathfrak{a}$  be determined by  $\langle H_\lambda, H \rangle = \lambda(H)$  for  $H \in \mathfrak{a}$ . For a Weyl chamber  $\mathfrak{a}^+ \subset \mathfrak{a}$  let  $\mathfrak{a}_+^*$  correspond to  $\mathfrak{a}^+$  under  $\lambda \rightarrow H_\lambda$ . Let  $\Sigma^+$  be the set of positive roots corresponding to  $\mathfrak{a}^+$ , let  $2\rho$  be their sum (with multiplicity) and  $\Sigma_0^+ = \Sigma^+ \cap \Sigma_0$ .

Let  $G = NAK$ ,  $\mathfrak{g} = \mathfrak{n} + \mathfrak{a} + \mathfrak{k}$  be the Iwasawa decompositions corresponding to our ordering. Let  $M$  be the centralizer of  $A$  in  $K$  and put  $B = K/M$ ,  $\overline{N} = \theta N$ . For  $H \in \mathfrak{a}$  let  $\partial(H)$  be the corresponding directional derivative. The map  $H \rightarrow \partial(H)$  extends to an isomorphism of  $S(\mathfrak{a})$  (resp.  $S^{\mathbf{C}} = \mathbf{C} \otimes S(\mathfrak{a})$ ) onto the algebra of all differential operators on  $\mathfrak{a}$  with constant real (respectively complex) coefficients. The map  $\lambda \rightarrow H_\lambda$  extends to an isomorphism of  $S^{\mathbf{C}}(\mathfrak{a}^*)$  onto  $S^{\mathbf{C}}(\mathfrak{a})$ .

Let  $I(\mathfrak{a})$  denote the algebra of  $W$ -invariants in  $S(\mathfrak{a})$  and  $I_+(\mathfrak{a})$  the set of  $q \in I(\mathfrak{a})$  without constant term. A polynomial  $p \in S(\mathfrak{a}^*)$  is said to be  $W$ -harmonic if  $\partial(q)p = 0$  for all  $q \in I_+(\mathfrak{a})$ . Let  $H(\mathfrak{a}^*)$  denote the space of  $W$ -harmonic polynomials and let  $H(\mathfrak{a}) \subset S(\mathfrak{a})$  correspond to  $H(\mathfrak{a}^*)$  under the extension of  $\lambda \rightarrow H_\lambda$ . It is well known that (see e.g. [GGA], Ch. III, §3)

$$(84) \quad S(\mathfrak{a}) = S(\mathfrak{a}) I_+(\mathfrak{a}) \oplus H(\mathfrak{a})$$

and

$$(85) \quad \dim H(\mathfrak{a}) = w.$$

For the adjoint action of  $K$  on  $\mathfrak{p}$  we have the analogous notions  $S(\mathfrak{p})$ ,  $S(\mathfrak{p}^*)$ ,  $I(\mathfrak{p})$ ,  $I(\mathfrak{p}^*)$ ,  $H(\mathfrak{p})$  and  $H(\mathfrak{p}^*)$ . The analog of (84) is also valid (loc. cit. III, §1).

As usual let  $\pi \in S(\mathfrak{a})$  be given by  $\pi = \prod_{\alpha \in \Sigma_0^+} H_\alpha$ , and  $\epsilon(\sigma)$  given by  $\pi^\sigma = \epsilon(\sigma)\pi$ ,  $\sigma \in W$ .

### Spectral Properties of $H(\mathfrak{a})$ .

With  $s_1, \dots, s_w$  running through  $W$  and  $p_1 = 1, p_2, \dots, p_w$  being a homogeneous basis of  $H(\mathfrak{a})$ , we consider the  $w \times w$  matrix

$E = (e_{ij})$  given by

$$(86) \quad e_{ij}(\lambda) = p_j(s_i\lambda) \quad \lambda \in \mathfrak{a}^*.$$

For simplicity we write  $S, I$  and  $H$  for  $S(\mathfrak{a}), I(\mathfrak{a})$  and  $H(\mathfrak{a})$ . We shall need the following properties of  $I$  and  $H$  from Chevalley [1955]. The proofs can also be found in [GGA], Ch. III, §3. The algebra  $I$  is generated by  $\ell$  algebraically independent homogeneous generators  $j_1, \dots, j_\ell$  of degrees  $d_1, \dots, d_\ell$ . The mapping  $j \otimes h \rightarrow jh$  extends to a linear bijection of  $I \otimes H$  onto  $S$ . With  $H^k$  denoting the space of homogeneous elements of  $H$  of degree  $k$  the Poincaré series of  $H$ ,  $\sum_{k \geq 0} (\dim H^k)t^k$ , has the form

$$(87) \quad P_H(t) = \prod_{1 \leq i \leq \ell} (1 + t + \dots + t^{d_i-1}).$$

Moreover,

$$(88) \quad w = d_1 \cdots d_\ell \quad \text{and} \quad \sum_i^\ell (d_i - 1) = r,$$

where  $r$  is the number of reflections in  $W$ . The sum  $\sum_1^w \deg(p_i)$  is clearly  $\sum_0^w k \dim H^k$ , which equals  $P'_H(1)$ . Taking the logarithmic derivative of  $P_H(t)$  at  $t = 1$  we get

$$(89) \quad P'_H(1)/P_H(1) = \frac{1}{2} \sum_1^\ell (d_i - 1) = \frac{r}{2}, \quad P_H(1) = w.$$

Given  $\lambda \in \mathfrak{a}^*$  put

$$S_\lambda = \{p \in S : p(\lambda) = 0\}, \quad I_\lambda = S \cap I.$$

The latter has codimension 1 in  $I$ . Also

$$S = IH \approx I \otimes H, \quad SI_\lambda = I_\lambda H \approx I_\lambda \otimes H$$

so we see from the above remark that  $SI_\lambda$  is an ideal of codimension  $w$  in  $S$  and that the  $p_i$  are linearly independent mod  $SI_\lambda$ . If  $\lambda$  belongs to the set  $\mathfrak{a}'$  of regular elements in  $\mathfrak{a}$  then the points  $s\lambda$  ( $S \in W$ ) are all different and  $S/SI_\lambda$  has  $w$  distinct homomorphisms into  $\mathbb{C}$ . Following a method kindly provided by David Vogan we shall prove that

$$(90) \quad \bigcap_{s \in W} S_{s\lambda} = SI_\lambda.$$

The inclusion  $\supset$  being obvious it suffices to prove that the left hand side has codimension at least  $w$  in  $S$ . For this we use the following result.

**Proposition.** *If  $R$  is a commutative ring and  $\alpha_1, \dots, \alpha_w$  distinct homomorphisms onto  $\mathbf{C}$  then the direct sum*

$$\alpha_1 \oplus \dots \oplus \alpha_w : R \rightarrow \mathbf{C} + \dots + \mathbf{C}$$

*is surjective.*

In order to deduce (90) we take  $\alpha_i$  as the homomorphism  $p \rightarrow p(s_i \lambda)$  ( $p \in S$ ). Then the kernel of the direct sum  $\oplus_i \alpha_i$  equals the left hand side of (90). Thus the proposition implies that the codimension is  $w$ .

For the proposition we observe the following easy result.

**Lemma.** *Suppose  $R$  is a commutative ring,  $M$  a maximal ideal and  $I \subset R$  an ideal not contained in  $M$ . Then the natural map of*

$$R/I \cap M \rightarrow R/I + R/M$$

*is an isomorphism.*

The proposition follows by induction taking  $M_i$  as the kernel of  $\alpha_i$  and taking  $I = m_1 \cap \dots \cap m_{i-1}$  and  $M = m_i$ .

**Lemma 5.18.**  $\text{Det}(E) = c\pi^{\frac{1}{2}w}$  where  $c \neq 0$  is a constant.

*Proof.* We first claim that  $\det(E(\lambda)) \neq 0$  if  $\lambda \in \mathfrak{a}'$ . If to the contrary  $\sum_{j=1}^w c_j p_j(s_i(\lambda)) = 0$  for some constants  $c_j$  (for all  $i$ ) then  $p = \sum c_j p_j$  vanishes on  $W \cdot \lambda$  so by (90)  $p \in SI_\lambda$ . But  $p \in H$  and  $H \cap SI_\lambda = 0$  so  $p = 0$ . If  $s_\alpha \in W$  is a reflection the map  $s \rightarrow ss_\alpha$  permutes the rows in  $E$ . Since  $p_j(ss_\alpha)(\lambda) = p_j(s\lambda)$  whenever  $H_\alpha(\lambda) = 0$  we see that the difference between the  $ss_\alpha$  row and the  $s$  row is divisible by  $H_\alpha$ . There being  $\frac{1}{2}w$  such pairs for each  $\alpha$ ,  $\det E$  is divisible by  $H_\alpha^{\frac{1}{2}w}$ . Since  $\alpha$  is arbitrary,  $\det E$  is divisible by  $\pi^{\frac{1}{2}w}$ . The same argument shows that each cofactor in  $E$  is divisible by  $\pi^{\frac{1}{2}w-1}$ .

On the other hand  $\deg(\det E) = \sum_j \deg p_j$ , which equals  $\sum_1^w j \dim H^j$  and by (88)–(89) this equals  $r \frac{w}{2}$ . This proves the lemma.

Consider now the quotient fields  $C(S)$  and  $C(I)$  and the bilinear form  $\langle \ , \ \rangle$  on  $C(S) \times C(S)$  with values in  $C(I)$  given by

$$(91) \quad \langle a, b \rangle = \sum_{\sigma \in W} a^\sigma b^\sigma.$$

Given  $(p_j) \subset H$  as above we can construct  $q^j \in C(S)$  explicitly satisfying

$$\langle p_j, q^i \rangle = \delta_{ij}$$

by making  $f_{k\ell}(\lambda) = q^k(s_\ell \lambda)$  the inverse of  $e_{\ell j}(\lambda) = p_j(s_\ell \lambda)$ . In fact

$$\delta_{kj} = \sum_{\ell} f_{k\ell} e_{\ell j} = \sum_{\ell} q^k(s_\ell \lambda) p_j(s_\ell \lambda) = \langle p_j, q^k \rangle.$$

**Lemma 5.19.** *Each  $q^j$  satisfies  $\pi q^j \in S$ .*

In fact by Lemma 5.18 and its proof  $q^j$  is the ratio of a cofactor divided by  $\det E$ . The lemma follows.

**Lemma 5.20.**

$$(i) \quad S = \bigoplus_{i=1}^w I p_i$$

$$(ii) \quad \frac{1}{\pi} S = \{x \in C(S(\mathfrak{a})) : \langle x, y \rangle \in I \text{ for } y \in S\}.$$

*Proof.* Part (i) is already proved. For (ii) let  $\check{I}$  denote the right hand side. Writing an  $\check{y} \in \check{I}$  in terms of the dual basis  $q^k$ ,  $\check{y} = \sum_k \frac{i_k}{j_k} q^k$ , we have  $\langle \check{y}, p_\ell \rangle = i_\ell / j_\ell$  so  $\check{y} \in \bigoplus_k I q^k$ . Thus by Lemma 5.19  $\pi \check{I} \subset S$ . On the other hand, let  $u \in S$ . Then for  $y \in S$ ,  $\langle \frac{u}{\pi}, y \rangle = \sum_\sigma \frac{u^\sigma}{\pi} \epsilon(\sigma) y^\sigma = \frac{1}{\pi} \sum_\sigma \epsilon(\sigma) (uy)^\sigma$  which belongs to  $S$  since sum is skew and thus divisible by  $\pi$ . Hence  $\frac{u}{\pi} \in \check{I}$ .

### 8. THE MULTITEMPORAL CAUCHY PROBLEM

The Iwasawa decomposition  $G = NAK$  gives a decomposition (Ch. III, §2, (18))

$$(92) \quad \mathbf{D}(G) = \mathbf{D}(A) \oplus (\mathfrak{n}\mathbf{D}(G) + \mathbf{D}(G)\mathfrak{k})$$

which provides a projection  $D \rightarrow q^D$  of  $\mathbf{D}(G)$  onto  $\mathbf{D}(A)$ . Restricting this to  $\mathbf{D}_K(G)$ , the algebra of  $K$ -right-invariant members of  $\mathbf{D}(G)$ , the mapping

$$D \rightarrow e^{-\rho} q^D \circ e^\rho$$

is a homomorphism of  $\mathbf{D}_K(G)$  onto  $I(\mathfrak{a})$  with kernel  $\mathbf{D}_K(G) \cap \mathbf{D}(G)\mathfrak{k}$  and this in turn gives an isomorphism  $\Gamma$  of  $\mathbf{D}(G/K)$  onto  $I(\mathfrak{a})$ . A member  $L$  of  $I(\mathfrak{a})$  will be regarded as a differential operator  $\partial(L)$  on  $\mathfrak{a}$ .

We now fix a real homogeneous basis  $p_1 = 1, p_2, \dots, p_n$  of  $H(\mathfrak{a})$ . For fixed  $f_i \in \mathcal{D}(X)$  ( $1 \leq i \leq w$ ) we now consider the “**multitemporal Cauchy problem**” for  $u \in \mathcal{E}(X \times \mathfrak{a})$ :

$$(93) \quad Du = \partial(\Gamma(D))u, \quad D \in \mathbf{D}(G/K)$$

$$(94) \quad (\partial(p_i)u)(x, 0) = f_i(x), \quad 1 \leq i \leq w$$

In (93)  $D$  operates on the first variable,  $\partial(\Gamma(D))$  on the second variable. As remarked by Schlichtkrull, the system (94) of initial data is equivalent to

$$(\partial(p)u)(x, 0) = F(p) \quad p \in H(\mathfrak{a}),$$

where  $F \in \text{Hom}_{\mathbf{R}}(H(\mathfrak{a}), \mathcal{D}(X))$ ,  $F(p_i) = f_i$ .

We shall now use the Fourier theory in Chapter III to study the system (93)–(94). We use the  $Q^\delta$  polynomial matrices given by III, §2 (22)–(22'). The representation  $\delta$  contragredient to  $\delta$  operates on the dual space  $V'_\delta = V_\delta$ . For  $v \in V_\delta$  let  $v' \in V'_\delta$  be determined by  $v'(u) = \langle u, v \rangle$  for  $u \in V_\delta$ . If  $C_\delta : V_\delta^M \rightarrow E_\delta$  is the linear transformation determined by  $C_\delta v'_j = \epsilon_j$  ( $1 \leq j \leq \ell(\delta)$ ) then

$$(95) \quad Q^\delta(\lambda)C_\delta v'_j = \sum_{i=1}^{\ell(s)} Q^\delta(\lambda)_{ij} v'_i.$$

We now prove a uniqueness theorem for the system (93)–(94).

**Proposition 5.21.** *Suppose  $u'$  and  $u''$  are two solutions to (93)–(94) and that for each  $H \in \mathfrak{a}$ , the functions  $x \rightarrow u'(x, H)$  and  $x \rightarrow u''(x, H)$  have compact support. Then  $u' = u''$ .*

*Proof.* Let  $u = u' - u''$ . Then  $u$  satisfies (93)–(94) with  $f_j \equiv 0$  for all  $j$ . Consider the Fourier transform  $\tilde{u}(\lambda, b; H)$  of  $u$  in the  $x$ -variable. Then if  $D \in \mathbf{D}(G/K)$  and  $D^*$  its adjoint,

$$(96) \quad \begin{aligned} \partial(\Gamma(D))_H(\tilde{u}(\lambda, b; H)) &= \Gamma(D^*)(-i\lambda)\tilde{u}(\lambda, b; H) \\ &= \Gamma(D)(i\lambda)\tilde{u}(\lambda, b; H) \end{aligned}$$

by Ch. III §1 (5). Also,

$$(97) \quad \{\partial(p_i)_H(u(x, H))\}_{H=0} \equiv 0.$$

Because of (84)–(85) each  $p \in S(\mathfrak{a})$  can be written  $p = \sum_{j=1}^w q_j p_j$  ( $q_j \in I(\mathfrak{a})$ ). Then using (93) and (96) we get

$$\begin{aligned} \partial(p)_H(\tilde{u}(\lambda, b; H)) &= \sum_j \partial(p_j)\partial(q_j)_H(\tilde{u}(\lambda, b; H)) \\ &= \sum_j \partial(p_j)_H(q_j(i\lambda)\tilde{u}(\lambda, b; H)) \end{aligned}$$

Using the inversion formula for  $u(x, H)$  and (97) we deduce

$$(98) \quad \{\partial(p)_H(u(x, H))\}_{H=0} = 0, \quad p \in S(\mathfrak{a}).$$

On the other hand, if  $\lambda \in \mathfrak{a}^*$  is regular (96) implies by a result of Steinberg and Harish-Chandra [GGA], III, Theorem 3.13 that for constants  $C_s$

$$\tilde{u}(\lambda, b; H) = \sum_{s \in W} C_s(\lambda, b) e^{is\lambda(H)}.$$

Thus the equation  $\{\partial(p)_H(\tilde{u}(\lambda, b; H))\}_{H=0} = 0$  for all  $p$  implies

$$\sum_{s \in W} p(is\lambda)C_s(\lambda, b) = 0, \quad p \in S(\mathfrak{a}).$$

For  $s_o \in W$  fixed we choose  $p$  such that  $p(is\lambda) = 0$  for all  $s \neq s_o$  in  $W$  and  $p(is_o\lambda) = 1$ . Then we get  $C_{s_o}(\lambda) = 0$ . Thus  $\tilde{u}(\lambda, b; H) = 0$  for all  $\lambda \in \mathfrak{a}^*$  regular and  $b \in B$  so by the inversion formula,  $u(x, H) \equiv 0$ , as desired.

Next we prove existence of the solution with an explicit formula.

**Theorem 5.22.** *Given  $f_1, \dots, f_w \in \mathcal{E}(X)$  the system*

$$(99) \quad Du = \partial(\Gamma(D))u \quad D \in \mathbf{D}(G/K)$$

*with initial data*

$$(100) \quad (\partial(p_j)u)(x, 0) = f_j(x) \quad (1 \leq j \leq w)$$

*has a solution given by the convolution*

$$(101) \quad u(x, H) = \sum_{j=1}^w (f_j \times S_H^j)(x)$$

*where for each  $j$ ,  $H \in \mathfrak{a}$ ,  $S_H^j$  is in  $\mathcal{E}'(X)$  and is for all  $b$  given by*

$$(102) \quad \sum_{\sigma \in W} q^j(i\sigma\lambda) e^{(i\sigma\lambda)(H)} = \int_X e^{(-i\lambda + \rho)(A(x, b))} dS_H^j(x).$$

*Remark.* Because of Prop. 5.21 we shall refer to (101) as *the* solution to the multitemporal Cauchy Problem.

Because of Lemma 5.19,  $q^j = h^j/\pi$  where  $h^j \in S(\mathfrak{a})$  so the left hand side of (102) equals

$$(103) \quad \frac{1}{\pi(i\lambda)} \sum_{\sigma} h^j(i\sigma\lambda) \epsilon(\sigma) e^{i\sigma\lambda(H)}$$

where  $\pi(\sigma\lambda) = \epsilon(\sigma)\pi(\lambda)$  and  $\epsilon(\sigma) = \pm 1$ . The sum in (103) is skew and is thus divisible by  $\pi(i\lambda)$  so (103) is holomorphic in  $\lambda$ . Also if  $\lambda = \xi + i\eta$  ( $\xi, \eta \in \mathfrak{a}^*$ ) we have

$$\left| e^{i\sigma\lambda(H)} \right| \leq e^{|H||\eta|}$$

so (103) is an entire function of exponential type of polynomial growth for  $\lambda$  real. By the Paley–Wiener theorem extended to  $\mathcal{E}'(X)$  (Ch. III, Corollary 5.9) there exist unique  $S_H^j \in \mathcal{E}'(X)$  satisfying (102) and

$$(104) \quad \text{supp}(S_H^j) \subset B_{|H|}(o).$$

Next we prove that  $S_H^j$  depends smoothly on  $H \in \mathfrak{a}$ , that is, for each  $f \in \mathcal{D}(X)$ ,  $S_H^j(f)$  is smooth in  $H$ . Since  $S_H^j$  is  $K$ -invariant,  $(f \times S_H^j)^\sim = \tilde{f}(S_H^j)^\sim$  (Ch. III, Lemma 1.4) so (102) implies for the lift  $\tilde{S}_J^j$  to  $G$

$$\int_G f(g^{-1} \cdot o) d\tilde{S}_H^j(g) = \int_{\mathfrak{a}^* \times B} \tilde{f}(\lambda, b) \sum_{\sigma} q^j(-i\sigma\lambda) e^{-i\sigma\lambda(H)} d\mu(\lambda, b)$$



so the smoothness in  $H$  follows.

Defining  $u(x, H)$  by (101) we shall now verify (99)–(100). Since  $D(f \times S) = f \times DS$  for  $f \in \mathcal{E}(X)$ ,  $S \in \mathcal{E}'(X)$  we consider  $DS_H^j$  and shall prove

$$(105) \quad DS_H^j = \partial(\Gamma(D))_H(S_H^j).$$

For this we observe that by (102)

$$(DS_H^j)\tilde{(\lambda, b)} = \Gamma(D^*)(-i\lambda)(S_H^j)\tilde{(\lambda, b)} = \Gamma(D)(i\lambda)(S_H^j)\tilde{(\lambda, b)}.$$

On the other hand, applying  $\partial(\Gamma(D))$  to (102) we see by the  $W$ -invariance of  $\Gamma(D)$  that

$$(\partial(\Gamma(D))_H(S_H^j))\tilde{(\lambda, b)} = \Gamma(D)(i\lambda)(S_H^j)\tilde{(\lambda, b)}.$$

These formulas imply (105), so (99) follows. Secondly, applying  $\partial(p_k)_H$  to (102) we get

$$\sum_{\sigma} q^j(i\sigma\lambda)p_k(i\sigma\lambda)e^{i\sigma\lambda(H)} = (\partial(p_k)_H(S_H^j))\tilde{(\lambda, b)}.$$

Putting here  $H = 0$  and using  $\langle q^j, p_k \rangle = \delta_{jk}$ , we obtain

$$(106) \quad \left\{ (\partial(p_k)_H(S_H^j))\tilde{(\lambda, b)} \right\}_{H=0} = \delta_{jk}$$

for all  $\lambda, b$ . By the injectivity of the Fourier transform this means

$$(107) \quad \left\{ \partial(p_k)_H(S_H^j) \right\}_{H=0} = \delta_{jk}\delta_o,$$

where  $\delta_o$  the delta distribution of  $X$  at  $o$ . Now (100) follows immediately.

**Remark.** Relations (105)–(107) are the analogs to (81)–(82) and one can thus consider the family

$$\{S_H^1, \dots, S_H^w\}$$

as the *fundamental solution* to (99)–(100).

**Example.**  $G/K$  of rank one. Let  $H_0$  be the vector in  $\mathfrak{a}^+$  of length 1. Then with  $p_1 = 1, p_2 = H_0$  we have  $q^1 = 1/2, q^2 = 1/2H_0$ . Also  $\Gamma(L) = H_0^2 - |\rho|^2$ . Writing  $v(x, t) = u(x, tH_0)$  equations (99) and (100) become

$$(L + |\rho|^2)v = \frac{\partial^2}{\partial t^2}v, \quad v(x, 0) = f_0(x), v_t(x, 0) = f_1(x)$$

Writing  $S_t^j$  instead of  $S_{tH_0}^j$  equation (102) becomes (since  $\lambda(H_0) = \pm|\lambda|$ )

$$\begin{aligned} \cos|\lambda|t &= \int_X e^{(-i\lambda+\rho)(A(x,b))} dS_t^1(x), \\ \frac{\sin|\lambda|t}{|\lambda|} &= \int_X e^{(-i\lambda+\rho)(A(x,b))} dS_t^2(x) \end{aligned}$$

in exact analogy with (80).

From Theorem 5.22 we can now deduce the following result stating informally that the speed of propagation is  $\leq 1$ .

**Corollary 5.23.** *If  $\text{supp}(f_j) \subset B_R(x_0)$  for  $1 \leq j \leq w$  then for each  $H \in \mathfrak{a}$ ,*

$$(108) \quad u(x, H) = 0 \quad \text{for } x \notin B_{R+|H|}(x_0).$$

By the group invariance we can take  $x_0 = o$ . Then the result is an immediate consequence of (101), (104) and the following simple lemma.

**Lemma 5.24** *Let  $f \in \mathcal{E}(X), S \in \mathcal{E}'(X)$ . Then*

$$(109) \quad \text{supp}(f) \subset \overline{B_r(o)}, \text{supp}(S) \subset \overline{B_s(o)} \Rightarrow \text{supp}(f \times S) \subset \overline{B_{r+s}(o)}.$$

*Proof.* By approximation we may assume  $S$  is a continuous function. If  $\tilde{S}$  denotes the lift of  $S$  to  $G$  ( $\tilde{S}(g) = S(g \cdot o)$ ) then

$$(f \times S)(g \cdot o) = \int_G f(gh^{-1} \cdot o) \tilde{S}(h) dh.$$

If this is  $\neq 0$  then for some  $h, d(gh^{-1} \cdot o, o) \leq r$  and  $d(h \cdot o, o) \leq s$  whence,

$$d(o, g \cdot o) = d(g^{-1} \cdot o, o) \leq d(g^{-1} \cdot o, h^{-1} \cdot o) + d(h^{-1} \cdot o, o) \leq r + s,$$

proving (109).

We can now reformulate the property in Cor. 5.23 as follows.

**Proposition 5.25** *Let  $f_j \in \mathcal{D}(X) (1 \leq j \leq w)$  and  $u(x, H)$  the corresponding solution to (99)–(100). Then if*

$$(110) \quad \overline{B_\rho(x_0)} \cap \bigcup_{j=1}^w \text{supp}(f_j) = \emptyset$$

*we have for each  $H \in \mathfrak{a}$ ,*

$$(111) \quad u(x, H) = 0 \quad \text{for } x \in B_{\rho-|H|}(x_0).$$

*Proof.* The compact set  $C = \bigcup_{j=1}^w \text{supp}(f_j)$  can be covered by finitely many balls  $B_\epsilon(x_i)$  ( $\epsilon$  fixed) such that

$$(112) \quad \overline{B_\rho(x_0)} \cap \overline{B_\epsilon(x_i)} = \emptyset \quad \text{for each } i.$$

Choose a corresponding partition of unity i.e.,  $\varphi_i \in \mathcal{D}(B_\epsilon(x_i)), \varphi_i \geq 0, \sum_i \varphi_i \leq 1$  with equality in a neighborhood of  $C$ . Then  $f_j = \sum_i f_j^i (1 \leq$

$j \leq w$ ) where  $f_j^i = \varphi_i f_j$  has support in  $B_\epsilon(x_i)$ . Let  $u^i(x, H)$  be the solution to (99)–(100) corresponding to  $f_j^i$  ( $1 \leq j \leq w$ ). By Cor. 5.23,

$$(113) \quad u^i(x, H) = 0 \quad \text{for } x \notin B_{\epsilon+|H|}(x_i).$$

However,  $B_{\epsilon+|H|}(x_i)$  is disjoint from  $\overline{B_{\rho-|H|}(x_0)}$ ; in fact a common point  $y$  would satisfy  $d(x_i, y) \leq \epsilon + |H|$ ,  $d(x_0, y) \leq \rho - |H|$  whence  $d(x_i, x_0) \leq \epsilon + \rho$ , contradicting (112). Thus by (113),  $u^i(x, H) = 0$  for  $x \in B_{\rho-|H|}(x_0)$  so (111) follows since  $u = \sum_i u^i$ .

We shall now derive some other formulas for the solution. If  $S \in \mathcal{E}'(X)$  is  $K$ -invariant and  $f \in \mathcal{D}(X)$  we have

$$(f \times S)^\sim = \tilde{f} \tilde{S}$$

so by (101) and the inversion formula for the Fourier transform,

$$u(x, H) = \frac{1}{w} \int_{\mathfrak{a}^*} \sum_{j, \sigma} q^j(i\sigma\lambda) e^{i\sigma\lambda(H)} (\mathbf{c}(\lambda)\mathbf{c}(-\lambda))^{-1} \int_B \tilde{f}_j(\lambda, b) e^{(i\lambda+\rho)(A(x,b))} db d\lambda.$$

Since  $(\mathbf{c}(\lambda)\mathbf{c}(-\lambda))^{-1}$  as well as the integral over  $B$  are  $W$ -invariant in  $\lambda$  (cf. Ch. III, Lemma 1.2), the terms in  $\sum_{\sigma}$  all have the same integral over  $\mathfrak{a}^*$ . Thus we can replace the average  $\frac{1}{w} \sum_{\sigma}$  by a single term. For each  $\sigma \in W$  we thus obtain,

$$(114) \quad u(x, H) = \int_{\mathfrak{a}^* \times B} \sum_k j_k^\sigma(\lambda) \tilde{f}_k(\lambda, b) e^{i\lambda(A(x,b)+\sigma H)} e^{\rho(A(x,b))} d\lambda db,$$

where

$$(115) \quad j_k(\lambda) = q^k(i\lambda)/\mathbf{c}(\lambda)\mathbf{c}(-\lambda)$$

and as usual,  $j_k^\sigma(\lambda) = j_k(\sigma^{-1}\lambda)$ . Consider now the Euclidean Fourier transform

$$F^*(\lambda) = \int_A e^{-i\lambda(\log a)} F(a) da$$

and let  $J_k$  denote the operator on  $\mathcal{S}(A)$  determined by

$$(116) \quad (J_k F)^*(\lambda) = j_k(\lambda) F^*(\lambda).$$

Since  $q^k(\lambda)\pi(\lambda)$  is a polynomial and  $(\pi(\lambda)\mathbf{c}(\lambda)\mathbf{c}(-\lambda))^{-1}$  has all its derivatives bounded by a polynomial,  $J_k$  does indeed map  $\mathcal{S}(A)$  into itself. Also, by a simple computation,

$$(J_k^\sigma F)^*(\lambda) = j_k^\sigma(\lambda) F^*(\lambda).$$

Consider now the Radon transform

$$\widehat{f}(\xi) = \int_{\xi} f(x) dm(x),$$

where  $\xi$  is a horocycle  $gNg^{-1} \cdot x$  and  $dm$  the measure on  $\xi$  derived from  $dn$ . More explicitly, if  $\xi_0 = N \cdot o$ ,

$$\widehat{f}(g \cdot \xi_0) = \int_N f(gn \cdot o) dn.$$

Since each  $\xi$  equals  $ka \cdot \xi_0$  where  $kM \in K/M$  and  $a \in A$  are unique we can write  $\widehat{f}(\xi)$  in the form  $\widehat{f}(kM, a)$ . Then  $J_k$  operates on functions on the horocycle space  $\Xi$ . Also by Ch. III, §5, (25)

$$(117) \quad \widetilde{f}(\lambda, b) = \int_A \widehat{f}(b, a) e^{(-i\lambda + \rho)(\log a)} da = (e^\rho \widehat{f})^*(\lambda, b).$$

Thus by (114)

$$u(x, H) = \int_{\mathfrak{a}^* \times B} \sum_k (J_k^\sigma(e^\rho \widehat{f}_k))^*(\lambda, b) e^{i\lambda(A(x, b) + \sigma H)} e^{\rho(A(x, b))} d\lambda db$$

so by the inversion formula for the Euclidean Fourier transform,

$$(118) \quad u(x, H) = \int_B \sum_k J_k^\sigma(e^\rho \widehat{f}_k)(b, \exp(A(x, b) + \sigma H)) e^{\rho(A(x, b))} db.$$

Here  $e^\rho$  is the function on  $\Xi$  given by  $e^\rho(kM, a) = e^{\rho(\log a)}$ . As in Ch. II, §3 we introduce the operators

$$\Lambda_{k, \sigma} = e^{-\rho} J_k^\sigma \circ e^\rho \quad (1 \leq k \leq w)$$

on the space  $\mathcal{S}_\rho(\Xi) = \{\varphi \in \mathcal{E}(\Xi) : e^\rho \varphi \in \mathcal{S}(\Xi)\}$ , where  $\mathcal{S}(\Xi) = \mathcal{S}(K/M \times A)$ . Then the space  $\mathcal{S}_\rho(\Xi)$  and the operators  $\Lambda_{k, \sigma}$  are  $G$ -invariant (Ch. II, Lemma 3.7.)

**Theorem 5.26** *Fix  $\sigma \in W$ . The solution to (99)-(100) can be written*

$$(119) \quad u(g \cdot o, H) = e^{\rho(\sigma H)} \int_{K/M} \sum_j (\Lambda_{j, \sigma} \widehat{f}_j)(gk \exp \sigma H \cdot \xi_0) dk_M.$$

*Proof.* If  $g = e$  the right hand side becomes

$$e^{\rho(\sigma H)} \int_B \sum_j (e^{-\rho} J_j^\sigma(e^\rho \widehat{f}_j))(b, \exp \sigma H) db = \int_B \sum_j (J_j^\sigma(e^\rho \widehat{f}_j))(b, \exp \sigma H) db,$$

which agrees with the right hand side of (118) for  $x = o$ . By the  $G$ -invariance of (99)-(100) in the  $x$ -variable the initial data  $f_j^{\tau(g^{-1})}$  generate the solution  $u^{\tau(g^{-1})}$ . Since the Radon transform as well as the operators  $\Lambda_{j,\sigma}$  commute with the  $G$ -action we deduce from the above

$$\begin{aligned} u(g \cdot o, H) &= u^{\tau(g^{-1})}(o, H) = e^{\rho(\sigma H)} \int_{K/M} \sum_j (\Lambda_{j,\sigma}(f_j^{\tau(g^{-1})})^\wedge)(k \exp \sigma H \cdot \xi_0) dk_M \\ &= e^{\rho(\sigma H)} \int_{K/M} \sum_j (\Lambda_{j,\sigma} \widehat{f_j})(gk \exp \sigma H \cdot \xi_0) dk_M \end{aligned}$$

as claimed.

Remark. For the case  $\text{rank}(X) = 1$ , Theorem 5.26 gives an alternative solution to the shifted wave equation in Theorem 5.7.

**Definition.** Given  $\sigma \in W$  the function

$$(S_\sigma f)(\xi) = \sum_j (\Lambda_{j,\sigma} \widehat{f_j})(\xi), \quad \xi \in \Xi$$

will be called the  $\sigma$ -source of the solution  $u$ .

**Corollary 5.27.** *The solution to (99)–(100) is*

$$(120) \quad u(x, H) = e^{\rho(\sigma H)} \int_B (S_\sigma f)(b, \exp(A(x, b) + \sigma H)) e^{2\rho(A(x, b))} db$$

and this formula is the symmetric space analog to (76).

From (120) it is easy to deduce the analog of Huygens' principle in §5, (6').

**Corollary 5.28.** *Assume  $G$  has all its Cartan subgroups conjugate. Then if  $\text{supp}(f_j) \subset B_R(o)$  ( $1 \leq j \leq w$ ) the solution  $u(x, H)$  has support in the region*

$$(121) \quad |H| - R \leq d(x, o) \leq |H| + R.$$

In fact, the assumption on  $G$  is equivalent to all  $\alpha \in \Sigma(\mathfrak{g}, \mathfrak{a})$  having even multiplicity ([DS], Ch. IX, §6). This in turn implies by the Harish-Chandra and Gindikin–Karpelevic formula for  $\mathfrak{c}(\lambda)$  that  $\mathfrak{c}(\lambda)^{-1}$  and even  $(\pi(\lambda)\mathfrak{c}(\lambda))^{-1}$  are polynomials. In addition,  $\pi(\lambda)q^j(i\lambda)$  is a polynomial so the operators  $\Lambda_{j,\sigma}$  are differential operators. Since

$$(122) \quad d(o, a \cdot o) \leq d(o, na \cdot o), \quad a \in A, n \in N,$$

we thus deduce from (120) that  $u(x, H) = 0$  if  $|H + A(x, b)| \geq R$  for  $b \in B$ . But if  $x = gK, b = kM, k^{-1}g = nak_0$ , (122) implies

$$(123) \quad |A(x, b)| = |A(k^{-1}g)| = |\log a| \leq d(o, na \cdot o) \leq d(o, x).$$

Thus if  $|H| > R + d(o, x)$  we have for  $b \in B$ ,

$$|H + A(x, b)| \geq |H| - |A(x, b)| \geq R$$

so

$$u(x, H) = 0 \text{ for } d(o, x) < |H| - R.$$

This, together with (108) implies the corollary.

**The Case  $G$  complex.** In this case we can give a simpler formula for the solution of the multitemporal Cauchy Problem. For  $h^j$  as in (103) we have  $\partial(h^j(e^{i\sigma\lambda})) = h^j(i\sigma\lambda)e^{i\sigma\lambda}$  so by (102)–(103)

$$(124) \quad \widetilde{S}_H^j(\lambda) = \frac{1}{\pi(i\lambda)} \partial(h_H^j) \left( \sum_{\sigma} \epsilon(\sigma) e^{i\sigma\lambda(H)} \right).$$

The density  $\delta$  is given by

$$\delta(\exp H) = \prod_{\alpha \in \Sigma^+} (e^\alpha - e^{-\alpha})^{m_\alpha}$$

and now the multiplicities  $m_\alpha$  equal 2 so we have (see e.g. [GGA], Ch. IV, §5)

$$\delta^{\frac{1}{2}}(a) = \sum \epsilon(\sigma) e^{\sigma\rho(\log a)}$$

and the spherical function  $\varphi_\lambda$  equals

$$\varphi_\lambda(a) = \frac{\pi(\rho)}{\pi(i\lambda)} \frac{\sum_{\sigma} \epsilon(\sigma) e^{i\sigma\lambda(\log a)}}{\delta^{\frac{1}{2}}(a)}.$$

Hence

$$(125) \quad \widetilde{S}_H^j(\lambda) = \frac{1}{\pi(\rho)} \partial(h^j)_H (\delta^{\frac{1}{2}}(\exp H) \varphi_\lambda(\exp H)).$$

**Theorem 5.29.** *For  $G$  complex the solution to (99)–(100) can be written*

$$(126) \quad u(x, H) = w \sum_{j=1}^w \partial(h^j)_H (\delta^{\frac{1}{2}}(\exp H) (M^{\exp H} f_j)(x)),$$

where  $M^x$  is the mean value operator (§2, (1)).

**Remark.** The analogy with the odd-dimensional flat case (3) is striking.

*Proof.* The mean value operator satisfies

$$M^h(e^{(i\lambda+\rho)(A(x,b))}) = \varphi_\lambda(h)e^{(i\lambda+\rho)(A(x,b))}$$

(Ch. II, Prop. 2.6) so by the inversion formula (Ch. III, Theorem 1.3)

$$(127) \quad (M^{\exp H} f_j) = \frac{1}{w} \int_{\mathfrak{a}^* \times B} \varphi_\lambda(\exp H) \tilde{f}_j(\lambda, b) e^{(i\lambda+\rho)(A(x,b))} d\mu(\lambda, b)$$

so by (125)

$$\begin{aligned} & \frac{1}{\pi(\rho)} \int_{\mathfrak{a}^*} \partial(h^j)_H (\delta^{\frac{1}{2}} \exp H) \varphi_\lambda(\exp H) \int_B e^{(i\lambda+\rho)(A(x,b))} \tilde{f}_j(\lambda, b) db d\mu(\lambda, b) \\ &= \int_{\mathfrak{a}^* \times B} \widetilde{S}_H^j(\lambda) e^{(i\lambda+\rho)A(x,b)} \tilde{f}_j(\lambda, b) d\mu(\lambda, b) = w(f_j \times S_H^j)(x). \end{aligned}$$

In view of (127) and the solution formula (101) this proves the theorem.

For a hyperbolic space  $\mathbf{H}^n$  ( $n$  odd) the shifted wave equation (52) is a special case of our rank-one example. The explicit solution as given by (50)–(51) is quite analogous to the solution (126). In the case  $n = 3$ , the isometry group  $G$  of  $\mathbf{H}^3$  is the complex group  $\mathbf{SL}(2, \mathbf{C})$  and in this case the two formulas will coincide.

### 9. INCOMING WAVES AND SUPPORTS.

In this section we shall prove the symmetric space analog of Cor. 5.16, relating incoming waves to support properties of the sources  $S_\sigma$ .

As usual, we put

$${}^+ \mathfrak{a} = \{H \in \mathfrak{a} : \langle H, H' \rangle > 0 \text{ for } H' \in \mathfrak{a}^+\}$$

and  $\mathfrak{a}^- = -\mathfrak{a}^+$ . Note that  ${}^+ \mathfrak{a} \cup \{0\}$  is closed ([GGA], Ch. IV, §6, (13)) and  $\mathfrak{a}^+ \subset {}^+ \mathfrak{a}$ .

**Definition.** Fix  $\sigma \in W$ . The wave  $u(x, H)$  in (120) is said to be  $\sigma$ -incoming if

$$(128) \quad u(x, H) = 0 \text{ for } d(o, x) < |H|, \quad \sigma H \in \mathfrak{a}^-.$$

To distinguish from balls in  $X$ , an open ball in  $\mathfrak{a}$  of radius  $r$  and center  $H$  is denoted  $V_r(H)$ .

We shall now, in analogy with Cor. 5.16, characterize  $\sigma$ -incoming waves in terms of their  $\sigma$ -sources.

**Theorem 5.30.** Fix  $\sigma \in W$  and  $f_j \in \mathcal{D}(X)$ ,  $(1 \leq j \leq w)$ . Then  $u$  is  $\sigma$ -incoming if and only if

$$(129) \quad \text{supp}(S_\sigma f)(b, \cdot) \subset \exp(\overline{+\mathfrak{a}}) \quad \text{for all } b \in B.$$

Naturally, formula (120) is basic to the proof. First we prove (129)  $\Rightarrow$  (128). Suppose  $H \in \mathfrak{a}$ ,  $\sigma H \in \mathfrak{a}^-$  and  $d(o, x) < |H|$ . Then by (123),  $A(x, b) \in V_{|H|}(0)$ . Now let  $H_0 \in +\mathfrak{a} + (-\sigma H)$ . Then  $\langle H_0 + \sigma H, -\sigma H \rangle \geq 0$  so

$$-|H_0||H| \leq \langle H_0, \sigma H \rangle \leq -|H|^2$$

so  $H_0 \notin V_{|H|}(0)$ . Hence  $A(x, b) \notin +\mathfrak{a} + (-\sigma H)$  so by (129) and (120),  $u(x, H) = 0$ . This proves (129)  $\Rightarrow$  (128).

For the converse we assume (129). For  $\varphi \in \mathcal{D}(\mathfrak{a})$  we consider the integral

$$(130) \quad v(x) = \int_{\mathfrak{a}} u(x, H)\varphi(H) dH.$$

Then by (114) we obtain

$$(131) \quad v(x) = \int_{\mathfrak{a}^* \times B} F_\lambda^\sigma(b)\psi^*(\lambda)e^{(i\lambda+\rho)(A(x,b))} d\lambda db,$$

where

$$\psi(H) = \varphi(-\sigma^{-1}H) \quad \text{and} \quad F_\lambda^\sigma(b) = \sum_k j_k^\sigma(\lambda)\tilde{f}_k(\lambda, b).$$

Note that

$$F_\lambda^\sigma(b) = h_\sigma(b, \cdot)^*(\lambda),$$

the Fourier transform (in  $A$ ) of the function

$$(132) \quad h_\sigma(b, a) = \sum_k J_k^\sigma(e^\rho \hat{f}_k)(b, a), \quad b \in B, a \in A.$$

Let  $\sigma^* \in W$  be the element interchanging  $\mathfrak{a}^+$  and  $\mathfrak{a}^-$  and put  $\tau = \sigma^{-1}\sigma^*$ . Then  $\sigma H \in \mathfrak{a}^- \iff H \in \tau\mathfrak{a}^+$  so by (128)

$$(133) \quad u(x, H) = 0 \quad \text{for } d(o, x) < |H|, H \in \tau\mathfrak{a}^+.$$

Now fix  $H_0 \in \tau\mathfrak{a}^+$  and consider the solution  $(x, H) \rightarrow u(x, H + H_0)$  to (93) whose initial data are

$$\{\partial(p_i)_H(u(x, H + H_0))\}_{H=0} = (\partial(p_i)u)(x, H_0).$$

We denote these by  $F_i(x)$  ( $1 \leq i \leq w$ ). Then (133) implies  $F_i(x) = 0$  for  $d(o, x) < |H_0|$ , i.e.,  $\text{supp}(F_i) \cap B_{|H_0|}(o) = \emptyset$ . Hence by Prop. 5.25,

$$(134) \quad u(x, H + H_0) = 0 \quad \text{for } x \in B_{|H_0|-|H|}(o).$$



Of course,  $H + H_0 \in V_{|H_0|-\epsilon}(H_0) \iff |H| < |H_0| - \epsilon$  which in turn implies

$$B_{|H_0|-|H|}(o) \supset B_\epsilon(o).$$

Thus (134) implies the following result.

**Lemma 5.31.** *Let  $\tau = \sigma^{-1}\sigma^*$ , fix  $H_0 \in \tau\mathfrak{a}^+$  and let  $\epsilon < |H_0|$ . Then*

$$u(x, H + H_0) = 0 \quad \text{for } d(o, x) < \epsilon \quad \text{if } H + H_0 \in V_{|H_0|-\epsilon}(H_0).$$

Now let

$$(135) \quad \mathfrak{a}_\epsilon(\tau) = \{H \in \mathfrak{a} : \langle H, H' \rangle > \epsilon |H'| \text{ for some } H' \in \tau\mathfrak{a}^+\}.$$

Then  $\mathfrak{a}_\epsilon(\tau) = \tau \cdot \mathfrak{a}_\epsilon(e)$ . Let  $H_0$  in Lemma 5.31 run along a fixed ray  $\{te_0 : t > 0, |e_0| = 1\}$  in  $\tau\mathfrak{a}^+$ . The balls  $V_{|H_0|-\epsilon}(H_0)$  will then fill up the half space  $\langle H, e_0 \rangle > \epsilon$ . As  $H_0$  varies in  $\tau\mathfrak{a}^+$  (and  $|H_0| > \epsilon$ ) these half spaces have union  $\mathfrak{a}_\epsilon(\tau)$ . Thus Lemma 5.31 implies the following result.

**Lemma 5.32.** *Let  $\tau = \sigma^{-1}\sigma^*$ ,  $\epsilon > 0$  and  $H \in \mathfrak{a}_\epsilon(\tau)$ . Then*

$$u(x, H) = 0 \quad \text{for } d(o, x) < \epsilon.$$

Now let  $\varphi$  in (130) satisfy  $\varphi \in \mathcal{D}(\mathfrak{a}_\epsilon(\tau))$ . Then

$$(136) \quad v(x) = 0 \quad \text{for } d(o, x) < \epsilon,$$

which by (131) amounts to

$$(137) \quad v(x) = \int_{\mathfrak{a}^*} \psi^*(\lambda) \mathcal{P}_\lambda(F_\lambda^\sigma)(x) d\lambda = 0,$$

where  $\mathcal{P}_\lambda$  is the Poisson transform defined by

$$(138) \quad (\mathcal{P}_\lambda F)(x) = \int_B e^{(i\lambda+\rho)(A(x,b))} F(b) db.$$

As in the proof in Ch. III, §2 of the injectivity criterion for the Poisson transform we exploit (137) by means of the equation

$$(139) \quad \{D_g(v(g \cdot o))\}_{g=e} = 0$$

for all  $D \in \mathbf{D}(G)$ , the algebra of left-invariant differential operators on  $G$ . The kernel in (138) is given by

$$e^{(i\lambda+\rho)(A(gK, kM))} = \zeta_{-\lambda}(k^{-1}g),$$

where  $\zeta_\lambda(g) = e^{(-i\lambda+\rho)(A(g))}$ . Fix  $\delta \in \widehat{K}_M$  (Ch. III, §2) and consider the operators  $\epsilon_j(v_i)$  there. By the proof of Theorem 2.6 in Ch. III, they satisfy

$$(140) \quad (\epsilon_j(v_i)\xi_{-\lambda})(k^{-1}) = \sum_{p=1}^{\ell(\delta)} \langle v_i, \delta(k)v_p \rangle Q^\delta(-\lambda)_{pj}.$$

Writing (131) and (137) in the form (with  $dk_M = db$ )

$$(141) \quad v(g \cdot o) = \int_{\mathfrak{a}^*} d\lambda \int_{K/M} \zeta_{-\lambda}(k^{-1}g)F(\lambda, kM) dk_M,$$

where  $F(\lambda, b) = h_\sigma(b, \cdot)^*(\lambda)\psi^*(\lambda)$ , we thus have

$$0 = \{\epsilon_j(v_i)_g(v(g \cdot o))\}_{g=e} = \int_{\mathfrak{a}^*} \sum_{p=1}^{\ell(\delta)} Q^\delta(-\lambda)_{pj} F_{pi}^\delta(\lambda) d\lambda.$$

Here

$$F_{pi}^\delta(\lambda) = \int_{K/M} \langle v_i, \delta(k)v_p \rangle F(\lambda, kM) dk_M.$$

Now  $F(\lambda, b) = (h_\sigma(b, \cdot) \times \psi)^*(\lambda)$ ,  $\times$  denoting convolution on  $A$  (and  $\mathfrak{a}$ ). Thus, defining

$$\zeta_{pi}(H) = \int_{K/M} h_\sigma(kM, \exp H) \langle v_i, \delta(k)v_p \rangle dk_M \quad (H \in \mathfrak{a})$$

we have for  $1 \leq i, j \leq \ell(\delta)$ ,

$$(142) \quad \sum_{p=1}^{\ell(\delta)} \int_{\mathfrak{a}^*} Q^\delta(-\lambda)_{pj} (\zeta_{pi} \times \psi)^*(\lambda) d\lambda = 0.$$

Define  $Q_{pj}^\delta \in S(\mathfrak{a})$  by  $Q_{pj}^\delta(\lambda) = Q^\delta(\lambda)_{pj}$  and  $P_{pj}^\delta \in S(\mathfrak{a})$  by  $P_{pj}^\delta(\lambda) = Q_{pj}^\delta(i\lambda)$ . Then, since  $P_{pj}^\delta(i\lambda) = Q^\delta(-\lambda)$ , (142) implies by the Fourier inversion formula,

$$\sum_{p=1}^{\ell(\delta)} [(\partial(P_{pj}^\delta)\zeta_{pi}) \times \psi](0) = 0,$$

so, if  $\check{\psi}(H) = \psi(-H)$ ,

$$(143) \quad \sum_{p=1}^{\ell(\delta)} \int_{\mathfrak{a}} (\partial(P_{pj}^\delta)\zeta_{pi})(H) \check{\psi}(H) dH = 0.$$

This holds whenever  $\varphi \in \mathcal{D}(\mathfrak{a}_\epsilon(\tau))$ , and  $\epsilon > 0$ .

Let  $\mathfrak{a}_0(\tau) = \lim_{\epsilon \rightarrow 0} \mathfrak{a}_\epsilon(\tau)$ . Then  $\mathfrak{a}_0(\tau) = \tau \mathfrak{a}_0(e)$  so the complement of  $\mathfrak{a}_0(\tau)$  is  $-\tau({}^+\mathfrak{a}) \cup \{0\} = \sigma^{-1}({}^+\mathfrak{a}) \cup \{0\}$ . Since  $\check{\psi}(H) = \varphi(\sigma^{-1}H)$ , (143) implies

$$(144) \quad \sum_{p=1}^{\ell(\delta)} \partial(P_{pj}^\delta) \zeta_{pi} = 0 \quad \text{on } \mathfrak{a} \setminus {}^+\mathfrak{a}.$$

The matrix  $P^\delta$  with entries  $P_{pj}^\delta \in S(\mathfrak{a})$  satisfies

$$(145) \quad \det(P^\delta) = P_c P^\delta,$$

where  $P_c$  is a matrix whose entries  $P_{pi}^c \in S(\mathfrak{a})$  are the cofactors of  $P^\delta$ . Applying  $\partial(P_{pj}^c)$  to (144) and summing on  $j$  we get by (145),

$$(\partial(\det P^\delta) \zeta_{pi})(H) = 0 \quad \text{for } H \in \mathfrak{a} \setminus {}^+\mathfrak{a}.$$

Now  $\det P^\delta(\lambda) = \det Q^\delta(i\lambda)$  which by Ch. III, Theorem 4.2 and Cor. 11.3 is a product of factors  $\langle \lambda, \alpha \rangle + c_{j,\alpha}$  where  $c_{j,\alpha} > 0$ . Here  $\alpha \in \sum_0^+$  and  $j$  runs through a certain finite set. With  $H_\alpha \in \mathfrak{a}$  given by  $\langle H_\alpha, H \rangle = \alpha(H)$  we thus have for  $F = \zeta_{pi}$ , which belongs to  $\mathcal{S}(\mathfrak{a})$ ,

$$(146) \quad \prod_{\alpha, j} (\partial(H_\alpha) + c_{j,\alpha}) F = 0 \quad \text{on } \mathfrak{a} \setminus {}^+\mathfrak{a}.$$

We pull out a single factor  $\partial(H_\alpha) + c$  ( $c > 0$ ) and consider the equation

$$(\partial(H_\alpha) + c)f = 0 \quad \text{in } \mathfrak{a} \setminus {}^+\mathfrak{a}, \quad f \in \mathcal{S}(\mathfrak{a}).$$

Then the function

$$(147) \quad g(H) = e^{c_0 \alpha(H)} f(H) \quad c_0 = c / \langle \alpha, \alpha \rangle$$

satisfies

$$(148) \quad \partial(H_\alpha)g = 0 \quad \text{in } \mathfrak{a} \setminus {}^+\mathfrak{a}.$$

Fix  $H_0 \in \mathfrak{a} \setminus \mathfrak{a}^+$ . Since  $-H_\alpha \in \mathfrak{a} \setminus {}^+\mathfrak{a}$  and since  $\mathfrak{a} \setminus {}^+\mathfrak{a}$  is star-shaped with respect to 0 we see that  $H_0 - tH_\alpha = t(\frac{H_0}{t} - H_\alpha)$  belongs to  $\mathfrak{a} \setminus {}^+\mathfrak{a}$  for  $t > 0$  sufficiently large. Then (148) implies that the function  $t \rightarrow g(H_0 - tH_\alpha)$  is a constant  $C$  for large  $t > 0$  so by (147)

$$(149) \quad f(H_0 - tH_\alpha) = C e^{-C_0 \alpha(H_0)} e^{C_0 \alpha(H_\alpha) t}.$$

This contradicts  $f \in \mathcal{S}(\mathfrak{a})$  unless  $f \equiv 0$  on  $\mathfrak{a} \setminus {}^+\mathfrak{a}$ . Iterating this argument for the factors in (146) we deduce

$$(150) \quad \zeta_{pi} \equiv 0 \quad \text{on } \mathfrak{a} \setminus {}^+\mathfrak{a}, \quad 1 \leq i, p \leq \ell(\delta).$$

Observe now that for each  $k_0 \in K$  assumption (128) is valid for the function  $u_{k_0} : (x, H) \rightarrow u(k_0 \cdot x, H)$ . Since  $A(k_0 \cdot x, b) = A(x, k_0^{-1} \cdot b)$  it is clear from (118) that replacing  $u$  by  $u_{k_0}$  amounts to replacing  $h_\sigma$  by the function  $(b, a) \rightarrow h_\sigma(k_0 \cdot b, a)$ . Thus (150) implies

$$(151) \quad \int_{K/M} h_\sigma(k_0 k M, \exp H) \langle v_i, \delta(k) v_p \rangle dk_M = 0, \quad H \in \mathfrak{a} \setminus {}^+ \mathfrak{a},$$

and here we put  $u = k_0 k, k = k_0^{-1} u$ . Since the vectors  $\delta(k_0) v_i (k_0 \in K, 1 \leq i \leq \ell(\delta))$  span  $V_\delta$ , (151) shows that the function  $uM \rightarrow h_\sigma(uM, \exp H)$  is orthogonal to all  $\langle v, \delta(k) v_p \rangle (1 \leq p \leq \ell(\delta), v \in V_\delta)$ . Since  $\delta \in \widehat{K}_M$  is arbitrary this proves

$$h_\sigma(b, \exp H) = 0 \quad \text{for } b \in B, H \in \mathfrak{a} \setminus {}^+ \mathfrak{a}$$

so (129) is proved.

### 10. ENERGY AND SPECTRAL REPRESENTATION.

We now restate the Cauchy Problem (93)–(94) into a vector formulation and use this to generalize the classical energy for the wave equation to the present context. We recall the definition.

Let  $p \in S(\mathfrak{a})$ . By Lemma 5.20 we have a matrix  $L_p = (L_p)_{ij}$  with entries in  $I(\mathfrak{a})$  such that

$$pp_j = \sum_1^w (L_p)_{ij} p_i, \quad 1 \leq j \leq w.$$

Let  $D_p = (D_p)_{ij}$  denote the corresponding matrix with entries in  $\mathbf{D}(G/K)$ , i.e.,  $\Gamma((D_p)_{ij}) = (L_p)_{ij}$ . Let  ${}^t D_p$  denote the transpose of  $D_p$ .

**Lemma 5.33.** *The Cauchy problem (93)–(94) is equivalent to the problem*

$$(152) \quad {}^t D_p \mu = \partial(p) \mu, \quad p \in S(\mathfrak{a}),$$

where the column vector  $\mu$  is a smooth map from  $X \times \mathfrak{a}$  to  $\mathbf{C}^w$  and

$$(153) \quad {}^t \mu(x, 0) = (f_1(x), \dots, f_w(x)).$$

More precisely, if  $u$  satisfies (93)–(94) then

$$(154) \quad {}^t \mu(x, H) = ((\partial(p_1)u)(x, H), \dots, (\partial(p_w)u)(x, H))$$

satisfies (152)–(153) above and if  $\mu$  satisfies (152)–(153) then  $u = \mu_1$  (first component) satisfies (93)–(94).

*Proof.* Since  $p_k = p_k p_1 = \sum_i (L_{p_k})_{i1} p_i$  we have  $(L_{p_k})_{i1} = \delta_{ik}$ . Thus, if  $\mu$  satisfies (152)–(153) then

$$\partial(p_k)\mu_1 = (\partial(p_k)\mu)_1 = ({}^t D_{p_k}\mu)_1 = \sum_i (D_{p_k})_{i1}\mu_i = \mu_k$$

so (94) holds. Also, if  $D \in \mathbf{D}(G/K)$ ,  $L_{\Gamma(D)}$  is a diagonal matrix,  ${}^t D_{\Gamma(D)} = D I_w$  so (152) implies  $D\mu_1 = \partial(\Gamma(D))\mu_1$ . Thus  $u = \mu_1$  satisfies (93).

On the other hand, assume  $u$  satisfies (93)–(94). Define  $\mu$  by (154). Then if  $p \in S(\mathfrak{a})$ ,

$$\begin{aligned} (\partial(p)\mu)_j &= \partial(p)\mu_j = \partial(p)\partial(p_j)u = \partial(pp_j)u \\ &= \partial\left(\sum_i (L_p)_{ij} p_i\right)(u) = \sum_i \partial(p_i)\partial(\Gamma(D_p)_{ij})u \\ &= \sum_i \partial(p_i)(D_p)_{ij}u = \sum_i (D_p)_{ij}\partial(p_i)u = ({}^t D_p\mu)_j \end{aligned}$$

so (152) holds; also (153) is obvious.

**Remark.** If (152) holds for  $p \in \mathfrak{a}$  it holds for all  $p \in S(\mathfrak{a})$ .

In fact, the map  $p \rightarrow L_p$  is an isomorphism of  $S(\mathfrak{a})$  onto an algebra of  $w \times w$  matrices with entries in  $I(\mathfrak{a})$ . Thus  $\partial(pq)\mu = {}^t D_{pq}\mu$ , justifying the remark.

Consider now the matrix  $A = (A_{ij})$  with entries in  $I(\mathfrak{a})$  given by

$$(155) \quad A_{ij} = (\pi q^j, \theta(\pi q^i)),$$

where  $(\cdot, \cdot)$  is defined in (91). Let  $\mathcal{A} = (\mathcal{A}_{ij})$  be the matrix with entries in  $\mathbf{D}(G/K)$  given by  $\Gamma(\mathcal{A}_{ij}) = A_{ij}$ .

**Definition.** Given  $u, v \in \mathcal{E}(X \times \mathfrak{a})$  the *energy form* is defined by

$$(156) \quad E(u, v; H) = \int_X ({}^t \mu \mathcal{A} \bar{\nu})(x, H) dx,$$

the integral assumed convergent. Here  ${}^t \mu, \nu$  are given by (133) from  $u$  and  $v$ , respectively.

It is a routine matter to verify that the function  ${}^t \mu \mathcal{A} \bar{\nu}$  is independent of the choice of basis of  $H(\mathfrak{a})$ . However, the homogeneity of each  $p_i$  plays a role in the proof of Theorem 5.34 and in Theorem 5.35.

In the rank-one example discussed before this reduces to

$$E(u, v; {}^t H_0) = \frac{\langle \pi, \pi \rangle}{2} \int_X (u_t \bar{v}_t - u(L + |\rho|^2)\bar{v}) dx$$

which up to the factor  $\langle \pi, \pi \rangle$  agrees with the usual formula in (70) and (72).

As used in the proof of (96) we have

$$\Gamma(D^*) = \theta\Gamma(D) \quad D \in \mathbf{D}(G/K).$$

Hence

$$\Gamma(\mathcal{A}_{ij}^*)(i\lambda) = \Gamma(\mathcal{A}_{ij})(-i\lambda) = A_{ij}(-i\lambda) = A_{ji}(i\lambda) = \Gamma(\mathcal{A}_{ji})(i\lambda)$$

so  $\mathcal{A}_{ij}^* = \mathcal{A}_{ji}$ . This implies the Hermitian symmetry,

$$E(u, v; H) = \overline{E(v, u; H)}.$$

While this relation is a simple formality, the positivity

$$E(u, u; H) \geq 0$$

which is a consequence of (162) below in connection with Theorem 5.34 is more subtle and seems to require the Fourier transform theory. Because of (162) we have a norm

$$\|F\| = \left\{ \int_X ({}^tF \mathcal{A}\overline{F})(x) dx \right\}^{1/2}, \quad {}^tF = (f_1, \dots, f_w)$$

on  $\mathcal{D}(X) \times \dots \times \mathcal{D}(X)$  ( $w$  times) which by the injectivity in Theorem 5.36 is strictly positive for  $F \neq 0$ .

**Theorem 5.34.** *Let  $u$  and  $v$  be the solutions to (93)–(94) with initial data  $(f_i) \subset \mathcal{D}(X)$ ,  $(g_j) \subset \mathcal{E}(X)$ . (Theorem 5.22). Then the energy form  $E(u, v; H)$  is constant in the time  $H \in \mathfrak{a}$ .*

*Proof.* Consider the matrix  $B = (B_{ij})$  given by

$$B_{ij} = (\pi q^i, \pi q^j) \in I(\mathfrak{a})$$

and  $J$  the diagonal matrix whose  $i^{\text{th}}$  entry is  $(-1)^{d_i}$  where  $d_i = \deg p_i = -\deg(q^i)$ . Then  $A_{ij} = (\pi q^j, \theta(\pi q^i)) = (-1)^{\deg(\pi)+d_i} B_{ij}$  so if  $\mathcal{B} = (\mathcal{B}_{ij})$  with  $\Gamma(\mathcal{B}_{ij}) = B_{ij}$ ,

$$(-1)^{\deg(\pi)} JB = A, \quad (-1)^{\deg(\pi)} J\mathcal{B} = \mathcal{A}.$$

Taking inner product with  $q^m$  and  $p_k$ , respectively, we derive the relations

$$\sum_i B_{ik} p_i = \pi^2 q^k, \quad p q^j = \sum_i (L_p)_{ji} q^i, \quad p \in S(\mathfrak{a}),$$

which imply quickly that the matrix  $L_p B$  is symmetric. Thus

$$L_p B = B {}^t L_p, \quad D_p \mathcal{B} = \mathcal{B} {}^t D_p.$$

In particular, take  $p = H \in \mathfrak{a}$  and observe that  $\partial(p) {}^t\mu = {}^t(\partial(p)\mu)$ . Then if  $u$  and  $v$  are solutions and  $\mu, \nu$  determined by (154) we have (by (152))

$$\begin{aligned} \partial(p)_H(E(u, v; H)) &= \int_X [({}^tD_p\mu)\mathcal{A}\bar{\nu} + {}^t\mu\mathcal{A}\partial(p)\bar{\nu}] (x, H) dx \\ &= (-1)^{\deg(\pi)} \int_X [({}^tD_p\mu)J\mathcal{B}\bar{\nu} + {}^t\mu J\mathcal{B}{}^tD_p\bar{\nu}] (x, H) dx. \end{aligned}$$

In the last term we replace  $\mathcal{B}{}^tD_p$  by  $D_p\mathcal{B}$  and put  $\xi = \mathcal{B}\bar{\nu}$ . The integrand is a sum of the terms

$$\begin{aligned} \sum_j ({}^tD_p\mu)_j (J\xi)_j &= \sum_{i,j} (D_p)_{i,j} \mu_i (-1)^{d_j} \xi_j \\ \sum_i ({}^t\mu)_i (JD_p\xi)_i &= \sum_{i,j} \mu_i (-1)^{d_i} (D_p)_{i,j} \xi_j \end{aligned}$$

so it would suffice to prove

$$\int_X [(D_p)_{i,j} \mu_i \xi_j + (-1)^{d_i - d_j} \mu_i (D_p)_{i,j} \xi_j] (x, H) dx = 0.$$

Since  $pp_j = \sum_i (L_p)_{ij} p_i$ ,  $(L_p)_{ij}$  is either 0 or homogeneous of degree  $1 + d_j - d_i$ . Also  $\Gamma(D^*) = \theta\Gamma(D)$  as used before, so

$$\Gamma((D_p)_{ij}^*) = \theta((L_p)_{ij}) = (-1)^{1+d_j-d_i} \Gamma((D_p)_{ij}),$$

whence

$$(D_p)_{ij}^* = (-1)^{1+d_j-d_i} (D_p)_{ij}.$$

Thus the last integral equals

$$\int_X [(D_p)_{i,j} \mu_i \xi_j - \mu_i (D_p)_{i,j}^* \xi_j] (x, H) = 0.$$

Hence  $\partial(p)_H(E(u, v; H)) = 0$  as claimed.

We consider now for each  $\sigma \in W$ ,  $\lambda \in \mathfrak{a}_C^*$ ,  $b \in B$  the particular solution  $\mu^\sigma(x, H; \lambda, b)$  of (152) given by (154) for

$$(157) \quad u^\sigma(x, H; \lambda, b) = e^{i\lambda(H) + (i\sigma\lambda + \rho)(A(x,b))}.$$

This is indeed a solution for

$$(Du^\sigma) = \Gamma(i\sigma\lambda)u^\sigma = \Gamma(i\lambda)u^\sigma, \quad D \in \mathbf{D}(G/K)$$

and

$$\partial(\Gamma(D))u^\sigma = \Gamma(i\lambda)u^\sigma.$$

We consider then the linear map (*the spectral representation*)

$$(158) \quad \mathcal{E}^\sigma : F(x) \rightarrow \int_X ({}^tF A \bar{\mu}^\sigma)(x, 0; \lambda, b) dx,$$

where  ${}^tF(x) = (f_1(x), \dots, f_w(x))$ ,  $f_i \in \mathcal{D}(X)$ . Thus  $\mathcal{E}^\sigma$  maps  $\mathcal{D}(X) \times \dots \times \mathcal{D}(X)$  into a function space on  $\mathfrak{a}^* \times B$ . The following result shows that  $\mathcal{E}^\sigma$  is intimately related to the Fourier transform on  $X$ .

**Theorem 5.35.** *For each  $\sigma \in W$*

$$\mathcal{E}^\sigma(F)(\lambda, b) = \pi(\lambda)^2 \sum_k q^k(i\lambda) \tilde{f}_k(\sigma\lambda, b).$$

*Proof.* For  $\lambda \in \mathfrak{a}^*$  (real) we have

$$\mu_j^\sigma(x, 0; \lambda, b)^- = p_j(-i\lambda) e^{(-i\sigma\lambda + \rho)(A(x, b))}$$

so

$$(159) \quad \mathcal{E}^\sigma(F) = \int_X \sum_{i, j} f_i(x) A_{ij}(-i\sigma\lambda) p_j(-i\lambda) e^{(-i\sigma\lambda + \rho)(A(x, b))} dx,$$

which by the  $W$ -invariance of  $A_{ij}$  equals

$$(160) \quad \sum_k \tilde{f}_k(\sigma\lambda, b) \sum_j A_{kj}(-i\lambda) p_j(-i\lambda).$$

Here the sum  $\sum_j$  equals

$$\begin{aligned} & \sum_j \left[ \sum_{\tau \in W} (\pi q^j)(-\tau i\lambda) (\pi q^k)(\tau i\lambda) \right] p_j(-i\lambda) \\ &= \sum_{\tau \in W} (\pi q^k)(\tau i\lambda) \sum_j (\pi q^j)(-\tau i\lambda) p_j(-i\lambda). \end{aligned}$$

Now if  $d_j = \text{degree } p_j$  then  $\text{degree } q^j = -d_j$  so the last  $\sum_j$  equals

$$(-1)^{\text{deg } \pi} \sum_j (\pi q^j)(\tau i\lambda) p_j(i\lambda).$$

Thus

$$\begin{aligned} \sum_j A_{kj}(-i\lambda) p_j(-i\lambda) &= (-1)^{\text{deg}(\pi)} \sum_j p_j(i\lambda) \sum_\sigma (\pi q^k)(\sigma i\lambda) (\pi q^j)(\sigma i\lambda) \\ &= (-1)^{\text{deg}(\pi)} \sum_j p_j(i\lambda) (\pi q^k, \pi q^j)(i\lambda). \end{aligned}$$



However,  $\sum_j p_j(\pi q^k, \pi q^j) = \pi^2 q^k$  as we see by taking inner product with  $q^m$  and noting that  $(\pi^2 q^k, q^m) = (\pi q^k, \pi q^m)$ . Putting these formulas together we see that expression (160) equals

$$\sum_k \tilde{f}_k(\sigma\lambda, b) (-1)^{\deg \pi} \pi(i\lambda)^2 q^k(i\lambda)$$

and this proves the theorem.

Given  $H_0 \in \mathfrak{a}$  let  $U_{H_0}$  denote the map of  $\mathcal{D}(X) \times \dots \times \mathcal{D}(X)$  into itself given by

$$(161) \quad U_{H_0} : \mu(x, 0) \rightarrow \mu(x, H_0),$$

$\mu$  satisfying (152). Since  $\partial(p)$  is invariant under translations of  $\mathfrak{a}$  it follows that for each  $H \in \mathfrak{a}$ ,  $U_{H_0}$  maps the function  $x \rightarrow \mu(x, H)$  into the function  $x \rightarrow \mu(x, H + H_0)$ .

In view of Theorem 5.35 the following result can be viewed as an enlargement of the Plancherel theorem for the Fourier transform on  $X$ . In fact it will be proved by the tools described in Ch. III, §5.

**Theorem 5.36.** *For each  $\sigma \in W$  the map  $\mathcal{E}^\sigma$  is an injective norm-preserving map of  $\mathcal{D}(X) \times \dots \times \mathcal{D}(X)$  onto a dense subspace of  $L^2(\mathfrak{a}^* \times B, d\lambda db / |\pi(\lambda)\mathbf{c}(\lambda)|^2)$ . Thus*

$$(162) \quad \int_{\mathfrak{a}^* \times B} \mathcal{E}^\sigma(F) \overline{\mathcal{E}^\sigma(F)}(\lambda, b) \frac{d\lambda db}{|\pi(\lambda)\mathbf{c}(\lambda)|^2} = \int_X ({}^t F A \overline{F})(x) dx.$$

Moreover,  $\mathcal{E}^\sigma$  intertwines  $U_{H_0}$  and the endomorphism

$$e(H_0) : \varphi(\lambda, b) \rightarrow e^{i\lambda(H_0)} \varphi(\lambda, b) \quad \text{of } L^2(\mathfrak{a}^* \times B, \frac{d\lambda db}{|\pi(\lambda)\mathbf{c}(\lambda)|^2}),$$

i.e.,

$$\mathcal{E}^\sigma U_{H_0} = e(H_0) \mathcal{E}^\sigma.$$

*Proof.* We separate the various statements.

1. *The norm identity (162).* We have by Theorem 5.35,

$$(163) \quad \mathcal{E}^\sigma(F) \overline{\mathcal{E}^\sigma(F)} = \pi(\lambda)^4 \sum_{k, \ell} q^k(i\lambda) q^\ell(-i\lambda) \tilde{f}_k(\sigma\lambda, b) \overline{\tilde{f}_\ell(\sigma\lambda, b)}$$

and by (12) in Ch. III, §1,

$$\int \tilde{f}_k(\sigma\lambda, b) \overline{\tilde{f}_\ell(\sigma\lambda, b)} db$$

is independent of  $\sigma$ . Thus the left hand side of (162) equals

$$\frac{1}{w} \int_{\mathfrak{a}^* \times B} \sum_{k, \ell} \pi(\lambda)^2 \tilde{f}_k(\lambda, b) \overline{\tilde{f}_\ell(\lambda, b)} \sum_{\sigma \in W} q^k(\sigma i \lambda) q^\ell(-\sigma i \lambda) \frac{d\lambda db}{|\mathbf{c}(\lambda)|^2}.$$

By writing  $\pi(\lambda)^2 = \pi(\sigma i \lambda) \pi(-\sigma i \lambda)$  our expression becomes

$$(164) \quad \frac{1}{w} \int_{\mathfrak{a}^* \times B} \sum_{k, \ell} \tilde{f}_k(\lambda, b) \overline{\tilde{f}_\ell(\lambda, b)} A_{\ell k}(i \lambda) \frac{d\lambda db}{|\mathbf{c}(\lambda)|^2}.$$

On the other hand if  $g \in \mathcal{D}(X)$  we have the inversion formula

$$g(x) = \frac{1}{w} \int_{\mathfrak{a}^* \times B} \tilde{g}(\lambda, b) e^{(i\lambda + \rho)(A(x, b))} |\mathbf{c}(\lambda)|^{-2} d\lambda db.$$

Also  $A_{\ell k}$  is real so  $A_{\ell k}$  has real coefficients and

$$(165) \quad (\mathcal{A}_{\ell k} g)^- = \mathcal{A}_{\ell k} \bar{g}, \quad \Gamma(\mathcal{A}_{\ell k})(i \lambda) = A_{\ell k}(i \lambda).$$

By the Plancherel formula and (165)

$$\int_X f(x) (\mathcal{A}_{k \ell} \bar{g})(x) dx = \frac{1}{w} \int_{\mathfrak{a}^* \times B} \tilde{f}(\lambda, b) (\mathcal{A}_{k \ell}(i \lambda) \tilde{g}(\lambda, b))^- |\mathbf{c}(\lambda)|^2 d\lambda db.$$

However, the definition (155) shows easily that for  $\lambda$  real,  $\overline{A_{k \ell}(i \lambda)} = A_{\ell k}(i \lambda)$  and we introduce this in the last integral. Taking now  $f = f_k$ ,  $g = f_\ell$  and summing on  $k, \ell$  we see that the right hand side of (162) equals (164). This proves (162).

2. *The injectivity.* For this suppose  $\mathcal{E}^\sigma(F) = 0$ . Then by Theorem 5.35

$$(166) \quad \sum_k q^k(i \lambda) \tilde{f}_k(\sigma \lambda, b) = 0.$$

Here we multiply by  $e^{(i\sigma\lambda + \rho)(A(x, b))}$  and integrate over  $B$ . Since this integral

$$(167) \quad \int_B \tilde{f}_k(\lambda, b) e^{(i\lambda + \rho)(A(x, b))} db$$

is  $W$ -invariant in  $\lambda$  we derive from (166) for each  $s \in W$ ,

$$\sum_k q^k(s i \lambda) \int_B \tilde{f}_k(\lambda, b) e^{(i\lambda + \rho)(A(x, b))} db = 0.$$

We multiply by  $p_\ell(si\lambda)$  and sum on  $s$ . This implies that (167) vanishes identically in  $x$ . By the injectivity (see Ch. III, Theorem 4.4) of the Poisson transform for  $\lambda \in \mathfrak{a}^*$  we have  $\tilde{f}_k(\lambda, b) \equiv 0$  whence by the injectivity of the Fourier transform,  $f_k \equiv 0$ .

The proof has the following consequence (Schlichtkrull).

**Corollary 5.37.**

$$f_k(x) = \int_{\mathfrak{a}^* \times B} p_k(i\lambda) \mathcal{E}^\sigma(F)(\lambda, b) e^{(i\sigma\lambda + \rho)(A(x, b))} \frac{d\lambda db}{\pi(\lambda)^2 |\mathbf{c}(\lambda)|^2}.$$

In fact, substitute for  $\mathcal{E}^\sigma(F)$  from Theorem 5.35 and use the  $W$ -invariance of (167). Again, replacing  $\lambda$  by  $s\lambda$  and summing on  $s$  the right hand side reduces to  $f_k(x)$  because of the inversion formula.

3. *The surjectivity.* Because of Theorem 5.35 this amounts to showing that if  $g(\lambda, b)$  satisfies

$$\int_{\mathfrak{a}^* \times B} \tilde{f}(\sigma\lambda, b) \pi(\lambda)^2 q^j(i\lambda) g(\lambda, b) \frac{d\lambda db}{|\pi(\lambda)\mathbf{c}(\lambda)|^2} = 0$$

for all  $f \in \mathcal{D}(X)$ , and all  $j$  ( $1 \leq j \leq w$ ) then  $g \equiv 0$ . Putting  $g^j(\lambda, b) = q^j(\sigma^{-1}i\lambda)g(\sigma^{-1}\lambda, b)$  we then have

$$(168) \quad \int_{\mathfrak{a}^* \times B} \tilde{f}(\lambda, b) g^j(\lambda, b) |\mathbf{c}(\lambda)|^{-2} d\lambda db = 0, \quad f \in \mathcal{D}(X), \quad 1 \leq j \leq w.$$

If we had this relation on  $\mathfrak{a}_+^* \times B$  the problem would be solved since by the Plancherel theorem the functions  $\tilde{f}(\lambda, b)$  are dense in  $L^2(\mathfrak{a}_+^* \times B)$ ,  $|\mathbf{c}(\lambda)|^{-2} d\lambda db$ . The fact that (168) holds for all  $j$  is decisive in passing from  $\mathfrak{a}_+^* \times B$  to  $\mathfrak{a}^* \times B$ .

Once again we invoke the representations  $\delta \in \widehat{K}_M$  and the results for the  $\delta$ -spherical transform described in Ch. III, §2. We specialize  $f$  in (168) to  $f \in \mathcal{D}_\delta(X)$  and put

$$(169) \quad g_\delta^j(\lambda) = \int_K g^j(\lambda, kM) \delta(k) dk, \quad g_\delta(\lambda) = \int_K g(\lambda, kM) \delta(k) dk.$$

If  $A, B(k) \in \text{Hom}(V_\delta, V_\delta)$  then

$$\int_K \text{Tr}(AB(k)) dk = \text{Tr}\left(A \int_K B(k) dk\right).$$

Using in addition  $\mathfrak{a}^* = \cup_{s \in W} s\mathfrak{a}_+^*$  as well as formula (48) (Ch. III, §5,) we derive from (168) the relation

$$(170) \quad \sum_{s \in W} \int_{\mathfrak{a}_+^*} \text{Tr}(\tilde{f}(s\lambda)g_\delta^j(s\lambda)) |\mathbf{c}(\lambda)|^{-2} d\lambda = 0.$$

Now according to Ch. III, Cor. 5.14 and (95) this section,

$$(171) \quad \tilde{f}(s\lambda) = Q^{\check{\delta}}(s\lambda)F_\delta(\lambda) = Q^{\check{\delta}}(s\lambda)C_\delta C_\delta^{-1}F_\delta(\lambda),$$

where  $F_\delta$  is a  $W$ -invariant holomorphic function on  $\mathfrak{a}_C^*$ , rapidly decreasing on  $\mathfrak{a}^*$  and with values in  $\text{Hom}(V_\delta, E_\delta)$ . We insert the expression (171) into (170). The map  $Q^{\check{\delta}}(s\lambda)C_\delta$  is a linear transformation of  $V_\delta^M$  into itself, but we extend it to  $V_\delta$  defining it 0 on the orthocomplement of  $V_\delta^M$  in  $V_\delta$ . Then  $g_\delta^j(s\lambda)$  can be shifted to the left in (170) and we obtain

$$(172) \quad \int_{\mathfrak{a}_+^*} \text{Tr}(G_j(\lambda)F_\delta^o(\lambda)) |\mathbf{c}(\lambda)|^{-2} d\lambda = 0,$$

where

$$(173) \quad G_j(\lambda) = \sum_{s \in W} g_\delta^j(s\lambda)Q^{\check{\delta}}(s\lambda)C_\delta, \quad \lambda \in \mathfrak{a}^*.$$

$$(174) \quad F_\delta^o(\lambda) = C_\delta^{-1}F_\delta(\lambda), \quad \lambda \in \mathfrak{a}^*.$$

Here  $F_\delta^o(\lambda)$  maps  $V_\delta$  into  $V_\delta^M$ . By the Paley–Wiener theorem for the  $\delta$ -spherical transform (Cor. 5.14, Ch. III) we can for each  $i$  and  $k$  ( $1 \leq i \leq \ell(\delta)$ ), ( $1 \leq k \leq d(\delta)$ ) take  $F_\delta^o(\lambda) = \varphi(\lambda)P^{ik}$  where  $\varphi(\lambda)$  is an arbitrary  $W$ -invariant holomorphic scalar-valued function of exponential type on  $\mathfrak{a}_C^*$ , rapidly decreasing on  $\mathfrak{a}^*$ , and  $P^{ik}(v_m) = \delta_{km}v_i$  ( $1 \leq m \leq d(\delta)$ ). Then

$$\text{Tr}(G_j(\lambda)F_\delta^o(\lambda)) = \varphi(\lambda)G_j(\lambda)_{ki}.$$

Thus by (172)

$$(175) \quad G_j(\lambda)_{ki} = 0 \quad 1 \leq i \leq \ell(\delta), \quad 1 \leq k \leq d(\delta),$$

for almost all  $\lambda \in \mathfrak{a}_+^*$ , hence by the  $W$ -invariance for almost all  $\lambda \in \mathfrak{a}^*$ . This means that  $G_j(\lambda)$  maps  $V_\delta^M$  into 0 and by definition it is 0 on  $(V_\delta^M)^\perp$ . Thus

$$(176) \quad G_j(\lambda) = 0 \text{ for almost all } \lambda \in \mathfrak{a}^*, \quad 1 \leq j \leq w.$$

We now have by (169)

$$(177) \quad g_\delta^j(s\lambda) = q^j(\sigma^{-1}si\lambda)g_\delta(\sigma^{-1}s\lambda), \quad 1 \leq j \leq w.$$

If  $a_{mn}(\lambda)$  is an arbitrary matrix entry in  $g_\delta(\sigma^{-1}\lambda)Q^\delta(\lambda)C_\delta^\zeta$  and  $s_1, \dots, s_w$  the elements in  $W$  then by (176)–(177),

$$(178) \quad \sum_{r=1}^w q^j(\sigma^{-1}s_r i\lambda)a_{mn}(s_r\lambda) = 0, \quad 1 \leq j \leq w.$$

By its definition the matrix  $\{q^j(s_r\nu)\}_{1 \leq j, r \leq w}$  is non-singular for regular  $\nu$  so (178) implies

$$g_\delta(\sigma^{-1}\lambda)Q^\delta(\lambda)C_\delta^\zeta = 0 \quad \text{for almost all } \lambda.$$

On  $V_\delta^M$ ,  $\det(Q^\delta(\lambda)C_\delta^\zeta) \neq 0$  so

$$(179) \quad g_\delta(\sigma^{-1}\lambda)v = 0 \quad \text{for } v \in V_\delta^M.$$

However, if  $u \in V_\delta$  we have for  $m \in M$

$$\begin{aligned} g_\delta(\lambda)u &= \int_K g(\lambda, kM)\delta(k)u \, dk \\ &= \int_K g(\lambda, kM)\delta(km)u \, dk = \int_M g(\lambda, kM)\delta(k) \, dk \int_M \delta(m)u \, dm \end{aligned}$$

which vanishes by (179). Thus  $g_\delta(\lambda) = 0$  for all  $\delta \in \widehat{K}_M$  whence  $g(\lambda, b) = 0$  for almost all  $\lambda \in \mathfrak{a}^*$  as desired.

4. *Intertwining.* Putting  ${}^tF = (f_1, \dots, f_w)$  we have by (161),

$$U_{H_0}(F)(x) = {}^t\mu(x, H_0) = (\partial(p_1)u, \dots, \partial(p_w)u)(x, H_0).$$

Taking Fourier transforms in (101) and then using (102) we get

$$\begin{aligned} (\partial(p_k)u)(\cdot, H_0)^\sim(\lambda, b) &= \partial(p_k)_{H_0} \left( \sum_j \tilde{f}_j(\lambda, b) \tilde{S}_{H_0}^j(\lambda, b) \right) \\ &= \sum_j \tilde{f}_j(\lambda, b) \sum_{\tau \in W} q^j(i\tau\lambda) p_k(i\tau\lambda) e^{i\tau\lambda(H_0)}. \end{aligned}$$

Using Theorem 5.35 we thus find

$$(180) \quad \mathcal{E}^\sigma U_{H_0}(F) = \pi(\lambda)^2 \sum_k q^k(i\lambda) \sum_j \tilde{f}_j(\sigma\lambda, b) \sum_{\tau \in W} q^j(i\tau\sigma\lambda) p_k(i\tau\sigma\lambda) e^{i\tau\sigma\lambda(H_0)}.$$

By construction the matrices

$$e_{ij} = p_j(s_i\lambda) \quad f_{kl} = q^k(s_\ell\lambda)$$

are inverses. Summing on the index  $k$  in the right hand side of (180) we see that

$$\sum_k q^k(i\lambda)p_k(i\tau\sigma\lambda) = 0 \text{ unless } \tau\sigma = e.$$

Thus (180) implies by  $\langle q^j, p_k \rangle = \delta_{jk}$ ,

$$\begin{aligned} \mathcal{E}^\sigma U_{H_0}(F)(\lambda, b) &= \pi(\lambda)^2 e^{i\lambda(H_0)} \sum_j \tilde{f}_j(\sigma\lambda, b) q^j(i\lambda) \\ &= e(H_0) \mathcal{E}^\sigma(F)(\lambda, b) \end{aligned}$$

as stated.

Because of Theorem 5.35 it is natural to describe the image of the spectral representation  $\mathcal{E}^\sigma$  by means of the Paley–Wiener theorem for the Fourier transform on  $X$ . Let  $\mathcal{H}(\mathfrak{a}^* \times B)$  be the space of functions satisfying (2) in Ch. III, §5 for some  $R$ .

**Theorem 5.38.** *The range of  $\mathcal{D}(X) \times \cdots \times \mathcal{D}(X)$  under  $\mathcal{E}^\sigma$  is the set of functions  $\varphi(\lambda, kM)$  on  $\mathfrak{a}_\mathbb{C}^* \times B$  satisfying*

- (i)  $\varphi \in \mathcal{H}(\mathfrak{a}^* \times B)\pi$ .
- (ii) For each  $\delta \in \widehat{K}_M$  the function

$$\lambda \rightarrow Q^\delta(\sigma\lambda)^{-1} \int_K \varphi(\lambda, kM) \delta(k^{-1}) dk$$

is holomorphic on  $\mathfrak{a}_\mathbb{C}^*$ .

*Proof.* We first prove the theorem with  $\sigma = e$ . Let  $\mathcal{E} = \mathcal{E}^e$ . In the proof we invoke the  $\delta$ -spherical transform (47) in Ch. III, §5. For  $f \in \mathcal{D}(X)$  let  $f^\delta \in \mathcal{D}(X, \text{Hom}(V_\delta, V_\delta^M))$  and  $\chi_s$  be as in Ch. III, §5. Then

$$(181) \quad \text{Tr}(f^\delta) = d(\delta)\chi_\delta * f \in \mathcal{D}_\delta(X)$$

and

$$(182) \quad \widetilde{\text{Tr}(f^\delta)}(\lambda) = (\tilde{f})^\delta(\lambda),$$

where on the left we have the  $\delta$ -spherical transform and on the right we have defined for each function  $\varphi$  on  $\mathfrak{a}_\mathbb{C}^* \times B$

$$\varphi^\delta(\lambda) = d(\delta) \int_K \varphi(\lambda, kM) \delta(k^{-1}) dk.$$

Let  $F = (f_1, \dots, f_w) \in \mathcal{D}(X) \times \dots \times \mathcal{D}(X)$  and put  $\varphi = \mathcal{E}(F)$ . Then (i) is obvious from Theorem 5.35. Moreover, (ii) is satisfied, since by Theorem 5.11, the function  $Q^{\check{\delta}}(\lambda)^{-1} \widetilde{\text{Tr}(f^\delta)}(\lambda)$  is holomorphic, and by (182) the latter function equals  $Q^{\check{\delta}}(\lambda)^{-1} \widetilde{f^\delta}(\lambda)$ .

Conversely, assume that a given function  $\varphi(\lambda, kM)$  satisfies (i) and (ii). Define

$$(183) \quad f_j(x) = \int_{\mathfrak{a}^* \times B} p_j(i\lambda) \varphi(\lambda, b) e^{(i\lambda + \rho)A(x,b)} \frac{d\lambda db}{\pi(\lambda)^2 |\mathbf{c}(\lambda)|^2}$$

for  $x \in X, j = 1, \dots, w$ . We claim that  $f_j \in \mathcal{D}(X)$  and  $\varphi = \mathcal{E}(F)$  where  $F = (f_1, \dots, f_w)$ .

We know that  $\lambda \mapsto \pi(\lambda)\mathbf{c}(\lambda)$  is nonzero on  $\mathfrak{a}^*$  and that its reciprocal satisfies a polynomial estimate. It then follows from the uniform, rapid decrease of  $\varphi(\cdot, b)$  that (183) defines a smooth function on  $X$ .

An easy computation shows that

$$f_j^\delta(x) = d(\delta) \int_K f_j(kx) \delta(k^{-1}) dk = \int_{\mathfrak{a}^*} p_j(i\lambda) \Phi_{\lambda, \delta}(x) \varphi^\delta(\lambda) \frac{d\lambda}{\pi(\lambda)^2 |\mathbf{c}(\lambda)|^2}$$

where  $\Phi_{\lambda, \delta}(x)$  is the generalized spherical function. The function  $\varphi^\delta(\lambda)$  is divisible by  $Q^{\check{\delta}}(\lambda)$  because of our assumption (ii), and we can then write

$$f_j^\delta(x) = \int_{\mathfrak{a}^*} p_j(i\lambda) \Phi_{\lambda, \delta}(x) Q^{\check{\delta}}(\lambda) Q^{\check{\delta}}(\lambda)^{-1} \varphi^\delta(\lambda) \frac{d\lambda}{\pi(\lambda)^2 |\mathbf{c}(\lambda)|^2}.$$

Averaging over  $W$ , and observing that  $\lambda \mapsto \Phi_{\lambda, \delta}(x) Q^{\check{\delta}}(\lambda)$  is  $W$ -invariant (cf. Ch. III, Theorem 5.15) we obtain

$$(184) \quad f_j^\delta(x) = \frac{1}{w} \int_{\mathfrak{a}^*} \Phi_{\lambda, \delta}(x) Q^{\check{\delta}}(\lambda) \left[ \sum_{s \in W} p_j(is\lambda) Q^{\check{\delta}}(s\lambda)^{-1} \varphi^\delta(s\lambda) \right] \frac{d\lambda}{\pi(\lambda)^2 |\mathbf{c}(\lambda)|^2}.$$

In order to prove that  $\text{Tr}(f_j^\delta)$  is compactly supported, it now suffices, by Ch. III, Cor. 5.14 and Theorem 5.15 to prove that the function

$$(185) \quad \lambda \mapsto \left[ \sum_{s \in W} p_j(is\lambda) Q^{\check{\delta}}(s\lambda)^{-1} \varphi^\delta(s\lambda) \right] \frac{1}{\pi(\lambda)^2}$$

belongs to the space  $\mathcal{J}^\delta(\mathfrak{a}^*)$  of  $W$ -invariant functions in  $\mathcal{H}(\mathfrak{a}_{\mathbf{c}}, \text{Hom}(V_\delta, V_\delta^M))$ . Clearly, (185) is  $W$ -invariant, and the expression in square brackets is holomorphic of exponential type, uniformly in  $\delta$ . It remains to be seen that it is divisible by  $\pi(\lambda)^2$ . We rewrite (185) as follows:

$$\left[ \sum_{s \in W} \epsilon(s) p_j(is\lambda) Q^{\check{\delta}}(s\lambda)^{-1} \frac{\varphi^\delta(s\lambda)}{\pi(s\lambda)} \right] \frac{1}{\pi(\lambda)}.$$

Here the expression in square brackets is holomorphic by (i) and (ii), since  $\pi(\lambda)$  and the determinant of  $Q^\delta(\lambda)$  have no common zeroes, cf. Ch. III, Theorem 4.2 and Cor. 11.3. Moreover, being skew, it is divisible by  $\pi(\lambda)$ . Hence  $\text{Tr}(f_j^\delta) \in \mathcal{D}(X)$  with support uniform in  $\delta$ , and so  $f_j = \sum_\delta \text{Tr}(f_j^\delta) \in \mathcal{D}(X)$ .

It remains to be seen that  $\mathcal{E}(F) = \varphi$ . It follows from the above that the  $\delta$ -spherical Fourier transform  $\widetilde{\text{Tr}(f_j^\delta)}$  of  $\text{Tr}(f_j^\delta)$  is given by  $Q^\delta(\lambda)$  times (185). By (182)  $\widetilde{\text{Tr}(f_j^\delta)} = \widetilde{f}_j^\delta$  and hence

$$\begin{aligned} \pi(\lambda)^2 \sum_j q^j(i\lambda) \widetilde{f}_j^\delta(\lambda) &= \sum_j q^j(i\lambda) Q^\delta(\lambda) \sum_{s \in W} p_j(is\lambda) Q^\delta(s\lambda)^{-1} \varphi^\delta(s\lambda) \\ &= \sum_{s \in W} \left( \sum_j q^j(i\lambda) p_j(is\lambda) \right) Q^\delta(\lambda) Q^\delta(s\lambda)^{-1} \varphi^\delta(s\lambda). \end{aligned}$$

Now  $\sum_j q^j(i\lambda) p_j(is\lambda) = 1$  for  $s = e$ , and otherwise it vanishes; hence we obtain

$$\pi(\lambda)^2 \sum_j q^j(i\lambda) \widetilde{f}_j^\delta(\lambda) = \varphi^\delta(\lambda).$$

Since this expression holds for all  $\delta$ , we conclude that

$$\pi(\lambda)^2 \sum_j q^j(i\lambda) \widetilde{f}_j(\lambda, b) = \varphi^\delta(\lambda, b)$$

for all  $b \in B$ , as claimed.

We now prove the theorem with arbitrary  $\sigma$ . We use the notation  $g^\sigma$ ,  $\varphi^\sigma$  for  $\lambda \mapsto g(\sigma^{-1}\lambda)$ ,  $(\lambda, b) \mapsto \varphi(\sigma^{-1}\lambda, b)$ . The proof will be based on the observation that conditions (i) and (ii) are independent of the basis  $p = (p_1, \dots, p_w)$  of  $H(a)$  used.

Let  $\mathcal{E}_p^\sigma$  denote the mapping  $\mathcal{E}^\sigma$  with the basis  $p$ . Thus

$$\mathcal{E}_p^e(\mathcal{D} \times \dots \times \mathcal{D}) = \mathcal{E}_{p'}^\sigma(\mathcal{D} \times \dots \times \mathcal{D})$$

if  $p' = (p'_1, \dots, p'_w)$  is any homogeneous basis of  $H(a)$  with  $p'_1 = 1$ . In particular,

$$(186) \quad \mathcal{E}_p^e(\mathcal{D} \times \dots \times \mathcal{D}) = \mathcal{E}_{p^\sigma}^e(\mathcal{D} \times \dots \times \mathcal{D}).$$

On the other hand

$$(\mathcal{E}_p^\sigma F)(\lambda, b) = \sum_k (q^k)^\sigma(i\sigma\lambda) \widetilde{f}_k(\sigma\lambda, b) \mathcal{E}_{p^\sigma}^e(\sigma\lambda, b)$$

so

$$(187) \quad \mathcal{E}_p^\sigma F = (\mathcal{E}_{p^\sigma}^e F)^{\sigma^{-1}}.$$

Using (186) and (187), the proof of (ii) is easily completed.



## 11. THE ANALOG OF THE FRIEDLANDER LIMIT THEOREM.

We shall now prove the analog for the multitemporal wave equation of Friedlander's Theorem 5.17. The proof is a reformulation of the proof in Phillips-Shahshahani [1993].

In the analogy with (116) let  $Q_k, J$  and  $\bar{J}$  denote the operators on  $\mathcal{S}(A)$  that under the Fourier transform on  $A$  correspond to multiplication by  $q^k(i\lambda)\pi(i\lambda)$ ,  $c(\lambda)^{-1}$  and  $c(-\lambda)^{-1}$ , respectively.

**Theorem 5.39.** *Let  $\sigma \in W$  and  $H, H_0 \in \mathfrak{a}^+$ , and put  $a_t = \exp(tH)$ . Then if  $k \in K$ ,*

$$\lim_{t \rightarrow +\infty} e^{t\rho(H)}(\partial(\pi)u)(ka_t \cdot o, -t\sigma H + \sigma H_0) = (\bar{J} \sum_k Q_k^{\sigma^{-1}}(e^\rho \hat{f}_k))(kM, H_0).$$

*Proof.* Because of (114)

$$(\partial(\pi)u)(x, H) = \int_{\mathfrak{a}^* \times B} F_\sigma(\lambda, b) e^{(i\lambda + \rho)(A(x, b))} e^{i\sigma\lambda(H)} \frac{d\lambda db}{|\mathbf{c}(\lambda)|^2},$$

where

$$F_\sigma(\lambda, b) = \sum_k \tilde{f}_k(\lambda, b) q^k(i\sigma\lambda) \pi(i\sigma\lambda).$$

Thus if  $\ell \in K$

$$(188) \quad \begin{aligned} & e^{t\rho(H)}(\partial(\pi)u)(\ell a_t \cdot o, -t\sigma H + \sigma H_0) \\ &= e^{t\rho(H)} \int_{\mathfrak{a}^* \times B} F_\sigma(\lambda, b) e^{(i\lambda + \rho)(A(a_t \cdot o, \ell^{-1}b))} e^{-it\lambda(H) + i\lambda(H_0)} \frac{d\lambda db}{|\mathbf{c}(\lambda)|^2}. \end{aligned}$$

By the  $K$ -invariance of  $db$ , we can replace

$$(189) \quad F_\sigma(\lambda, b) \text{ by } F_\sigma(\lambda, \ell \cdot b), \quad A(a_t \cdot o, \ell^{-1} \cdot b) \text{ by } A(a_t \cdot o, b).$$

In the Iwasawa decomposition  $G = KAN$ , we write as before

$$(190) \quad g = k(g) \exp H(g) n(g), \quad H(g) \in \mathfrak{a}.$$

Then  $A(gK, kM) = -H(g^{-1}k)$ . If  $b = k'M$  and we define  $k \in K$  by  $k' = k(a_t k)$  then  $kM \rightarrow k'M$  is a diffeomorphism and

$$(191) \quad dk' = e^{-2\rho(H(a_t k))} dk.$$

Since  $a_t k = k' \exp H(a_t k) n$  and  $A$  normalizes  $N$  we have

$$(192) \quad H(a_{-t} k') = -H(a_t k).$$

After the substitution (189) and  $A(a_t \cdot o, k'M) = -H(a_{-t}k')$  relations (191)–(192) reduce the right hand side of (188) to

$$(193) \quad e^{t\rho(H)} \int_{\mathfrak{a}^*} e^{-i\lambda(tH-H_0)} \int_B F_\sigma(\lambda, \ell k(a_t k)M) e^{(i\lambda-\rho)(H(a_t k))} \frac{d\lambda dk_M}{|\mathbf{c}(\lambda)|^2}.$$

Now we use the map  $\bar{n} \rightarrow k(\bar{n})M$  of  $\bar{N}$  into  $K/M$  under which  $dk_M = e^{-2\rho(H(\bar{n}))} d\bar{n}$  (see e.g. [GGA], Ch. I, §5) to transfer (194) to  $\bar{N}$ . Since  $A$  normalizes  $N$  we have

$$a\bar{n} = ak(\bar{n}) \exp H(\bar{n})n_1 = k(ak(\bar{n})) \exp H(ak(\bar{n})) \exp H(\bar{n})n_2$$

so we have Harish–Chandra’s relations

$$(194) \quad k(a\bar{n}) = k(ak(\bar{n})) = k(\bar{n}^a), \quad H(a\bar{n}) = H(ak(\bar{n})) + H(\bar{n})$$

$$(195) \quad H(a\bar{n}) - \log a = H(\bar{n}^a),$$

where  $\bar{n}^a = a\bar{n}a^{-1}$ . Thus (193) becomes

$$(196) \quad \int_{\mathfrak{a}^*} e^{-i\lambda(tH-H_0)} \int_{\bar{N}} F_\sigma(\lambda, \ell k(\bar{n}^{a_t})M) L(a_t, \bar{n}) \frac{d\lambda d\bar{n}}{\mathbf{c}(\lambda)\mathbf{c}(-\lambda)},$$

where the kernel  $L(a, \bar{n})$  equals

$$(197) \quad L(a, \bar{n}) = e^{(i\lambda-\rho)(H(\bar{n}^a))} e^{-(i\lambda+\rho)(H(\bar{n}))}.$$

Let  $\lambda = \xi + i\eta$  ( $\xi, \eta \in \mathfrak{a}^*$ ) and  $-\eta \in \mathfrak{a}_+^*$  so small that  $(\mathbf{c}(\lambda)\mathbf{c}(-\lambda))^{-1}$  is holomorphic in  $\mathfrak{a}^* + it\eta$  ( $|t| \leq 1$ ). Let  $\epsilon > 0$  be such that

$$0 < \epsilon < 1, \quad \rho + \epsilon\eta \in \mathfrak{a}_+^*.$$

Then by [GGA], Ch. IV, §6, (41),

$$|L(a_t, \bar{n})| \leq e^{(\epsilon\eta-\rho)(H(\bar{n}))},$$

which is integrable over  $\bar{N}$ . Since  $F_\sigma$  is holomorphic in  $\lambda$  and for a fixed  $\eta$  is rapidly decreasing in  $\xi$  the integral (196) can be shifted to  $\mathfrak{a}^* + i\eta$ . Letting then  $t \rightarrow +\infty$  we get the expression

$$\begin{aligned} & \int_{\mathfrak{a}^* + i\eta} e^{i\lambda(H_0)} (\mathbf{c}(\lambda)\mathbf{c}(-\lambda))^{-1} F_\sigma(\lambda, \ell M) \left( \int_{\bar{N}} e^{-(i\lambda+\rho)(H(\bar{n}))} d\bar{n} d\lambda \right) \\ & = \int_{\mathfrak{a}^* + i\eta} e^{i\lambda(H_0)} F_\sigma(\lambda, \ell M) (\mathbf{c}(-\lambda))^{-1} d\lambda. \end{aligned}$$

Here we can shift the integration back and take  $\eta = 0$ . Using the definition of  $\bar{J}$ ,  $Q_k$  and relation (116) the theorem follows.

To summarize we now compare the flat wave equation in  $\mathbf{R}^n$  with the multitemporal wave equation by the following table.

## Remarkable Analogies

<p><b>The Wave Equation in <math>\mathbf{R}^n</math></b></p> $Lu = \frac{\partial^2 u}{\partial t^2},$ $u(x, 0) = f_0(x), u_t(x, 0) = f_1(x)$	<p><b>Multitemporal Wave Equation in <math>G/K</math></b></p> $Du = \partial(\Gamma(D))u,$ $\partial(p_i)u(x, 0) = f_i(x)$
<p><b>Radon Transform Solution</b></p> $u(x, t) = \int_{\mathbf{S}^{n-1}} (Sf)(w, \langle x, w \rangle + t) dw$ $Sf = c(\partial^{n-1} \widehat{f}_0 + \partial^{n-2} \widehat{f}_1)$	<p><b>Horocycle Transform Solutions</b></p> $u(x, H) = e^{\rho(H)}.$ $\int_B (S_\sigma f)(b, \exp(A(x, b) + \sigma H)) e^{2\rho(A(x, b))} db$ $(S_\sigma f)(\xi) = \sum_{j=1}^w (\Lambda_{j, \sigma} \widehat{f}_j)(\xi).$
<p><b>Fourier Transform Solution</b></p> $u(x, t) = (f_0 * T'_t)(x, t) + (f_1 * T_t)(x, t)$ $\frac{\sin  \lambda  t}{ \lambda } = \int_{\mathbf{R}^n} e^{-i\langle x, \lambda \rangle} dT_t(x)$	<p><b>Fourier Transform Solution</b></p> $u(x, H) = \sum_{j=1}^w (f_j * S_H^j)(x)$ $\sum_{\sigma} q^j(i\sigma\lambda) e^{i\sigma\lambda(H)} =$ $\int_X e^{(-i\lambda + \rho)(A(x, b))} dS_H^j(x)$
<p><b>Mean value solution, <math>n</math> odd</b></p> $u(x, t) = c_n \left[ \frac{\partial}{\partial t} (I_t f_0)(x) + (I_t f_1)(x) \right]$ $(I_t f)(x) = \left( \frac{\partial}{\partial(r^2)} \right)^{\frac{n-3}{2}} (r^{n-2} (M^r f)(x))$	<p><b>Mean value solution, <math>G</math> complex</b></p> $u(x, H) =$ $w \sum_{j=1}^w \partial(h^j)_H [\delta^{\frac{1}{2}}(\exp H) (M^{\exp H} f_j)(x)]$
<p><math>u(x, t)</math> incoming if</p> $u(x, t) = 0 \text{ in }  x  < -t \quad t < 0$ <p>For <math>\mathbf{R}^{2n+1}</math>:</p> <p><math>u(x, t)</math> is incoming if and only if</p> $(\partial \widehat{f}_0 + \widehat{f}_1)(w, s) = 0 \text{ for } s < 0, \text{ all } w$	<p><math>u(x, H)</math> <math>\sigma</math>-incoming if</p> $u(x, H) = 0 \text{ for } d(0, x) <  H , \quad \sigma H \in \mathfrak{a}^-$ <p>For <math>\mathbf{H}^3</math>:</p> <p><math>u(x, H)</math> is <math>e</math>-incoming if and only if</p> $[H_0(e^\rho \widehat{f}_0) + c_n e^\rho \widehat{f}_1](b, H) = 0 \text{ for}$ <p>all <math>b, H \in \mathfrak{a}^-</math>.</p>