# on regular reflection of a shock wave from a rigid wall* 

## V.m. TESHUKOV

The three-dimensional problem of reflection of a shock wave of arbitrary front shape from a fixed rigid wall is considered. Existence of a piecewise analytic solution of the problem defining the initial stage of regular reflection is proved. Expansion of solution behind the reflected wave front in converging power series in the neighborhood of the incident wave (of the moving along the wall intersection line of the incident wave front and the rigid wall) is obtained for prolonged instants of time. It is shown that such expansions generally occur only then, when the length of the trace velocity vector relative to gas behind the reflected wave projected on the plane normal to the trace exceeds the speed of sound. A bibliography of publications dealing with the problem of shock wave reflection can be found in $/ 1,2 /$.

1. Statement of the problem. The motion of an inviscid non-heat-conducting gas is considered in the region of $\varphi(x) \geqslant 0$ with surface $\Gamma_{1}: \varphi(x)=0(\nabla \varphi \neq 0)$ assumed to be a rigid impervious wall at which the gas velocity vector $u$ satisfies the relation

$$
\begin{equation*}
\mathbf{u} \nabla_{\varphi}=0, \quad \mathbf{x}=(x, y, z) \in \Gamma_{1} \tag{1.1}
\end{equation*}
$$

Let the piecewise-analytic solution of gasdynamics equations

$$
\begin{equation*}
\rho \frac{d \mathbf{u}}{d t}+\nabla p=0, \quad \frac{d p}{d t}+\rho c^{2} \operatorname{div} \mathbf{u}=0, \quad \frac{d s}{d t}=0, \quad \rho=\psi(p, s) \tag{1.2}
\end{equation*}
$$

defining shock wave propagation toward a rigid wall when $t \leqslant 0$, be known. Here $p$ is the pressure, $\rho$ is the density, $s$ is the entropy, $c$ is the speed of sound, and $\psi(p, s)$ is an analytic function that specifies the equation of state of a normal gas $/ 3,4 /$. This means that ahead and behind the shock wave front $u(x, t), p(x, t), s(x, t)$ are analytic functions of their arguments, and the shock wave surface $\Gamma_{2}$ is an analytic hypersurface in the four-dimensional space $\mathbf{x}, t$. The Hugoniot relations

$$
\begin{equation*}
\left[\rho v_{n}\right]=0, \quad\left[p+\rho v_{n}^{2}\right]=0, \quad\left[\varepsilon+p \rho^{-1}+{ }^{1} / 2 v_{n}^{2}\right]=0, \quad\left[u_{\sigma}\right]=0 \tag{1.3}
\end{equation*}
$$

and the condition of entropy increase are satisfied on $\Gamma_{2}$. Here $[f]$ denotes the jump of quantity $f$ at transition through the discontinuity, $v_{n}=D_{n}-u_{n}, u_{n}$ and $D_{n}$ are the velocities of gas and of shock wave front, respectively, in the direction of the normal $\mathbf{n}$ to the front, $\mathbf{u}_{\boldsymbol{o}}$ is the tangent velocity component, and $\varepsilon$ is the specific internal energy. The solution ahead of the front satisfies condition (1.4). At instant $t=0$ the shock wave front reaches the wall touching it at point $Q$. Further motion of gas has to be defined.

When $t>0$ the solution structure changes, a shock wave reflected from the wall makes its appearance. The region of determination of the generalized gas motion in the space $x$, $t$ consists of three subregions, viz., region $\Omega_{1}$ bounded by $\Gamma_{1}$, the incident wave surface $\Gamma_{2}$ and the plane $t=0$; region $\Omega_{2}$ bounded by $\Gamma_{2}$, the reflected wave surface $\Gamma_{3}$ and the plane
$t=0$; and region $\Omega_{3}$ bounded by $\Gamma_{1}$ and $\Gamma_{3}$ (Fig.l illustrates the two-dimensional case). Solution in $\Omega_{1}$ and $\Omega_{2}$ is obtained idenpendently by solving the problem of an arbitrary discontinuity on the curvilinear surface $\gamma_{20}$ (the intersection of $\Gamma_{2}$ and the plane $t=0$ ), and the incident wave surface $\Gamma_{2}$ is determined for $t>0 / 5 /$. In what follows the solution in
$\Omega_{1}$ and $\Omega_{2}$, and surface $\Gamma_{2}$ are assumed known. We have to construct the solution of Eqs. (1.2) in $\Omega_{3}$ which satisfies conditions (1.1) and (1.3) on $\Gamma_{1}, \Gamma_{3}$, and at the same time determine the reflected wave surface $\Gamma_{3}$.
2. Relationships at the shock wave. In the case of regular reflection which obtairs in the initial stage the unknown surface $\Gamma_{3}$ must pass through the known two-dimensional surface $\gamma_{0}$, the locus of $\Gamma_{1}$ and $\Gamma_{2}$ intersection points (Fig.1). Let surface $\gamma_{0}$ be parametrically defined: $t=t_{0}(\beta, \gamma), \mathbf{x}=\mathbf{x}_{0}(\beta, \gamma)$ with $t_{0}$ and $\mathbf{x}_{0}$ being analytic functions of parameters $\beta, \gamma$, and $x_{0 \beta} \neq 0, x_{0 \gamma} \neq 0, x_{0 \beta} \times x_{0 \gamma} \neq 0$. We introduce in region $\Omega_{3}$ new coordinates
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$\tau, \alpha, \beta, \gamma$ such that $\tau-0$ corresponds to the shock wave surface, and $\alpha-0$ to the rigid wall. The respective substitution of variables conforms to $t=t(\tau, \alpha, \beta, \gamma)=\tau-\alpha+t_{0}(\beta, \gamma), \mathbf{x}=$ $\mathbf{x}(\tau, \alpha, \beta, \gamma)$. Function $\mathbf{y}=\mathbf{x}(0, \alpha, \beta, \gamma)$ is determined by the solution of the Cauchy problem

$$
\mathbf{y}_{\alpha}=\mathbf{H},\left.\quad \mathbf{y}\right|_{\alpha=0}=\mathbf{x}_{0}(\boldsymbol{\beta}, \boldsymbol{\gamma})
$$

Function $\mathbf{H}$ is selected so that equations $t=\alpha \dot{\alpha} t_{0}(\beta, \gamma), \mathbf{x}=\mathrm{x}(0, \alpha, \beta, \gamma)$ parametrically define the shock wave surface $\Gamma_{3}$.

Let the equation of $\Gamma_{3}$ be of the form $f_{1}(x, l)=0$. Then $\varphi_{1}(x(0, \alpha, \beta, \gamma), t(0, \alpha, \beta, \gamma))=0$ identically with respect to $\alpha, \beta, \gamma$. Differentiation of this equality with respect to $\alpha, \beta, \gamma$ yields the formulas

$$
\left(\varphi_{1}\right)_{t} \because x_{\alpha} \Gamma_{\varphi_{1}}=0, \quad t_{\beta}\left(\varphi_{1}\right)_{t}+\mathbf{x}_{\beta} \Gamma_{\varphi_{1}}=0, \quad t_{\gamma}\left(\varphi_{1}\right)_{t}+\mathbf{x}_{\gamma} \nabla \varphi_{1}=0
$$

Since the shock wave front normal to $n$ and the quantity $D_{n}$ are linked with $q_{1}$ by the relations $\mathbf{n} \cdots \nabla_{f_{1}}\left|\nabla_{\Upsilon_{1}}\right|^{1}, D_{n}=:-\left(\varphi_{1}\right)_{t}\left|\nabla_{\varphi_{1}}\right|^{-1}$, we have the following equalities:

$$
\mathbf{x}_{\beta} \mathbf{n}=t_{n} D_{u}, \quad \mathbf{x}_{\gamma} \mathbf{n}=t_{\gamma} D_{n}, \quad|\mathbf{n}|=1, \quad \mathbf{x}_{\gamma} \mathbf{n}=D_{n}
$$

of which the first three formulas enable us to determine $\mathbf{n}$ in the form

$$
\begin{equation*}
\mathbf{n}=\left\{\left(|\mathbf{m}|^{2}-D_{n}^{2}|\mathbf{k}|^{2}\right)^{1 / 2} \mathbf{m}+D_{n}(\mathbf{k} \times \mathbf{m})\right\}|\mathbf{m}|^{-2}, \quad \mathbf{m}=\mathbf{x}_{\beta} \times \mathbf{x}_{\gamma}, \quad \mathbf{k}=t_{\beta} \mathbf{x}_{\gamma}-l_{\gamma} \mathbf{x} \beta \tag{2.1}
\end{equation*}
$$

where $\mathbf{k}$ is a vector tangent to the shock wave front by virtue of (2.1).
Function $\mathbf{H}$ must, therefore, satisfy the relation $\mathbf{H n}=D_{n}$, where the normal $\mathbf{n}$ is defined by formula (2.1). Specific selection will be made in Sect.3. Writing the preceding formulas in the form

$$
\left(\mathbf{x}_{\beta}-t_{\beta} \mathbf{u}\right) \mathbf{n}=t_{\beta}\left(t_{n}-u_{n}\right), \quad\left(\mathbf{x}_{\gamma}-t_{\gamma} \mathbf{u}\right) \mathbf{n}=t_{\gamma}\left(D_{n}-u_{n}\right), \quad|\mathbf{n}|=\mathbf{1}
$$

we obtain for $n$ another formula

$$
\begin{align*}
& \left.\mathbf{n}=\left\{|\mathbf{q}|^{2}-v_{n}^{2}|\mathbf{k}| j^{2}\right)^{2} \mathbf{q}+v_{n}(\mathbf{k} \times \mathbf{q})\right\}|\mathbf{q}|^{\mathbf{2}}  \tag{2.2}\\
& \mathbf{q}=\left(\mathbf{x}_{\beta}-t_{\beta} \mathbf{1}\right) \times\left(\mathbf{x}_{\gamma}-t_{\gamma} \mathbf{u}\right)
\end{align*}
$$

which is valid when $|\mathbf{q}|>\left|v_{n} \mathbf{k}\right|$ and the inequality $|\mathbf{m}| \geqslant\left|D_{n} \mathbf{k}\right|$ is satisfied in (2.1). If function $\mathbf{x}(0, \alpha, \beta, \gamma)$ is known, $\mathbf{x}(\tau, \alpha, \beta, \gamma)$ is determined similarly by solving the Canchy problem

$$
\mathbf{x}_{\tau}=\mathbf{G},\left.\quad \mathbf{x}\right|_{\tau=0}=\mathbf{x}(0, \alpha, \beta, \gamma)
$$

Function $G$ is selected so that the contact characteristic $/ 3 /$ that passes at $\tau=0$ through the cross section $\alpha=$ const of surface $\Gamma_{3}$ is parametrically specified by formulas $t=\tau+\alpha+$ $t_{0}(\beta, \gamma), \mathbf{x}=\mathbf{x}(\tau, \alpha, \beta, \gamma)$ for fixed $\alpha$. If the equation of that characteristic is of the form $\varphi_{2}(x, t)=0$, then as previously we obtain the equalities

$$
\left(\varphi_{2}\right)_{t}+\mathbf{u} \nabla \varphi_{2}=0, \quad\left(\varphi_{2}\right)_{t}+\mathbf{x}_{\tau} \nabla \varphi_{2}=0, \quad t_{\beta}\left(\varphi_{2}\right)_{t}+\mathbf{x}_{\beta} \nabla \varphi_{2}=0, \quad t_{\gamma}\left(\varphi_{2}\right)_{t}+\mathbf{x}_{\gamma} \nabla \varphi_{\varphi_{2}}=0
$$

and their corollaries

$$
\left(\mathbf{x}_{\beta}-t_{\beta} \mathbf{u}\right) \mathbf{n}_{2}=0, \quad\left(\mathbf{x}_{\gamma}-t_{\mathbf{\gamma}} \mathbf{u}\right) \mathbf{n}_{2}=0, \quad\left(\mathbf{x}_{\tau}-\mathbf{u}\right) \mathbf{n}_{\mathbf{2}}=0, \quad\left|\mathbf{n}_{\mathbf{2}}\right|=1
$$

Hence $\quad \mathbf{n}_{\mathbf{2}}=\mathbf{q}|\mathbf{q}|^{-1}$ and function $\mathbf{G}$ must satisfy the relation $(\mathbf{G}-\mathbf{u}) \mathbf{q}=0$. The specific


Fig. 1


Fig. 2 selection of $G$ is effected in Sect. 3. Since $\mathbf{k q}=\mathbf{k n}=0$, vector $k$ at points $\gamma_{0}$ is directed along the tangent to the line of intersection of the incident wave front and the wall, i.e. to the trace which the incident front leaves on the wall at instant $t=$ const. A three-dimensional picture of this at instant $t=$ const appears in Fig. 2 ( $\gamma_{2 t}$, $\gamma_{3 t}, \gamma_{0 t}$ are cross sections of $\Gamma_{2}$, $\Gamma_{3}, \gamma_{0}$ at $t=$ const).

For analyzing the relationships at the shock wave it is convenient to use the vector $w$

$$
\mathbf{w}=(\mathbf{q} \times \mathbf{k})|\mathbf{k}|^{-\mathbf{2}}=\mathbf{u}-(\mathbf{u} \mathbf{k})|\mathbf{k}|^{-2} \mathbf{k}-(\mathbf{k} \times \mathbf{m})|\mathbf{k}|^{-2}
$$

which at points $\gamma_{0}$ represents the difference of projections of $\mathbf{u}$ on the plane normal to $\mathbf{k}$ and the trace motion velocity on the wall in a direction perpendicular to the trace. Since by virtue of (2.2) $w_{n}=\mathbf{w n}=-v_{n}$, hence on $\Gamma_{3}$

$$
\left[\rho w_{n}\right]=0,\left[p+\rho w_{n}^{2}\right]=0,\left[\varepsilon+p \rho^{-1}+1 / 2 w_{n}^{2}\right]=0,\left[w_{\sigma}\right]=0
$$

From the corollaries of relations at discontinuity

$$
|\mathbf{w}|^{2}=\left|\mathbf{w}_{1}\right|^{2}-\left(p-p_{1}\right)\left(v_{1}+v\right),\left(\mathbf{w}-\mathbf{w}_{1}\right)^{2}=\left(p-p_{1}\right) \times\left(v_{1}-v\right)
$$

where $v=\rho^{-1}$, and the quantities ahead of the front are denoted by subscript 1 and those behind it have no subscripts. For the angle $\theta$ of turn of $w$ in the $\Pi$ plane orthogonal to $k$ we have at transition through the discontinuity $\left(\cos \theta=\left(w w_{1}\right)|\mathbf{w}|^{-1} \times\left|w_{1}\right|^{-1}\right)$ the relation

$$
\begin{gather*}
\sin ^{2} \theta=\left(p-p_{1}\right)\left\{v_{1}-v-v_{1}^{2}\left|\mathbf{w}_{1}\right|^{-2}\left(p-p_{1}\right)\right\} \times\left\{\left|\mathbf{w}_{1}\right|^{2}-\left(p-p_{1}\right)\left(v_{1}+v\right)\right\}^{-1}  \tag{2.3}\\
v=g\left(p, p_{1}, v_{1}\right)
\end{gather*}
$$

where $\theta$ is the $p$-polar.
Function $g$ defines the Hugoniot curve. In the $\theta, p$ plane curve (2.3) is determinate when $\left(v_{1}-v\right)\left|\mathbf{w}_{1}\right|^{2} v_{1}^{-2} \geqslant p-p_{1} \geqslant 0$. Along that curve $|\sin \theta|<1$ and $\theta$ vanishes when equality is reached in one of inequalities. Consequently $\theta$ reaches its maximum value $\theta_{*}$ (limit angle of turn) at some point $\quad p=p_{*}$.

In what follows, we assume that the equations of state of the gas are such that the righthand side of (2.3) is monotonic when $p<p_{*}$ and $p>p_{*}$ and that there exists a $p_{0}\left(p_{1}<p_{0}<\right.$ $\left.p_{*}\right)$ such that along the curve $|\mathbf{w}|^{2}=\left|\mathbf{w}_{1}\right|^{2}-\left(p-p_{1}\right)\left(v_{1}+v\right)<c^{2}$ when $p>p_{0}$ and $|\mathbf{w}|^{2}>c^{2}$ when $p<p_{0}$. A polytropic gas satisfies these conditions.

Using the identity

$$
\sin ^{2} \theta=\left|\mathbf{q} \times \mathbf{q}_{1}\right|^{2}|\mathbf{q}|^{-2}\left|\mathbf{q}_{1}\right|^{-2}, \quad \mathbf{q}=\{(\mathbf{q} \mathbf{m}) \mathbf{m}+(\mathbf{u m})(\mathbf{k} \times \mathbf{m})\}|\mathbf{m}|^{-2}
$$

we obtain from (2.3) the relation

$$
\begin{equation*}
\frac{h-h_{1}}{\left\{\left(1+\left.h^{2}|\mathbf{k}|\right|^{2}\left(1+h_{2}^{2}|\mathbf{k}|=\mid\right)^{1 / 2}\right.\right.}=\left\{\frac{\left.\left(p-p_{1}\right)\left(p_{1}-v\right)-|\mathbf{k}| v_{1}^{2}\left|\mathbf{q}_{1}\right|^{-2}\left(p-p_{1}\right)^{2}\right\}^{1 / 2}}{\left|\boldsymbol{q}_{1}\right|^{2}-\left(p-\rho_{1}\right)\left(v_{1}+v\right)|\mathbf{k}|^{2}}\right\}^{2} \tag{2.4}
\end{equation*}
$$

At point $Q$, where the shock wave front touches the wall, ( $t_{\mathrm{p}}-t_{\gamma}-0, \mathrm{k}-0$ ), and (2.4) becomes the equation of the $\left(p, u_{n}\right)$-diagram of shock waves, which is determinate and monotonic for all values of $p\left(p_{*} \rightarrow \infty\right.$ as $\left.|\mathbf{k}| \rightarrow 0\right)$. In the region of determination (2.4)

$$
|\mathbf{q}|^{2}-v_{n}^{2}|\mathbf{k}|^{2}=|\mathbf{k}|^{2}\left(|w|^{2}-i_{n}^{2}\right)=|\mathbf{k}|^{2}\left(\left|\mathbf{w}_{1}\right|^{2}-\left(p-p_{1}\right)\left(v_{1}-v\right)^{-2} v_{1}^{2}\right) \geqslant 0
$$

From (1.3) with allowance for the equalities

$$
\left[\mathbf{x}_{\beta} \mathbf{u}\right]=\left(\mathbf{x}_{\beta} \mathbf{n}\right)\left[u_{n}\right]=t_{\beta} D_{n}\left[u_{n}\right], \quad\left[\mathbf{x}_{\gamma} \mathbf{u}\right]=t_{\gamma} D_{n}\left[u_{n}\right]
$$

we obtain two corollaries

$$
\begin{align*}
& {[a]=|b|=0 \quad\left(a=\left(\mathbf{u x}_{\beta}\right)-t_{\beta}\left(\varepsilon+p v \cdot 1_{2}|u|^{2}\right)\right.}  \tag{2.5}\\
& \left.b=\left(u x_{\gamma}\right)-t_{y}\left(\varepsilon+p v+1 / 2|u|^{2}\right)\right)
\end{align*}
$$

When $\mathbf{x}_{\beta}, \mathbf{x}_{\gamma}, t_{\beta}, t_{\gamma}$ at some point of surface $\Gamma_{3}$, the parameters of gas ahead of the front, and the quantity $h\left(h<h_{*}\right)$ behind it are known, it is possible to determine all parameters of gas behind the front, vector n and $D_{n}$.

In conformity with (2.1), at point $Q$ where $t_{\beta}=t_{\gamma}=0, \mathbf{n}=\boldsymbol{m}|\mathbf{m}|^{-1}$ we have $h=u_{n}|\mathbf{m}|^{-1}$. Equation (2.4) uniquely determines $p$ in terms of $h$, and subsequently all remaining parameters of gas $/ 3,4 /$. At points $\Gamma_{3}$, where $t_{p}{ }^{2}+t_{\nu}^{2} \neq 0$, Eq. (2.4) has for a given $h$ two solutions $p_{-}, p_{+}\left(p_{-}\right.$ $<p_{*}<p_{+}$) (an analog of "weak" and "strong" shocks in the theory of steady motion). But, when $t_{\rho}^{2}+t_{\gamma}{ }^{2} \rightarrow 0$, only $p_{-}$approaches $p$ which is uniquely defined at point $Q\left(p_{+} \rightarrow \infty\right.$ as $\left.|k| \rightarrow 0\right)$. Hence on considerations of continuity of solution behind the front we uniquely select the "weak" shock. Using the Hugoniot adiabatic curve and the known value of $p$, we determine $r$, and using known $p, v, p_{1}, v_{1}$ we obtain $i_{n}$. The following identities apply:

$$
\begin{gathered}
(\mathbf{u m})=h|\mathbf{q}||\mathbf{m}|\left(1+\left.h^{2}\left|\mathbf{k} \|^{-1 / 2}, \quad\right| \mathbf{q}\right|^{2}=|\mathbf{m}|^{2}+2 \mathbf{k}\left(b \mathbf{x}_{\beta}-a \mathbf{x}_{\gamma}\right)-2|\mathbf{k}|^{\mu}(\varepsilon+p v)-\left(t_{\beta} b-t_{\gamma^{n}}\right)^{z}\right. \\
\mathbf{u}=\|(\mathbf{u m}) \mathbf{m}+\left(u \mathbf{x}_{\gamma}-b \mathbf{x}_{\beta}\right) \times \mathbf{m}+\left(v+p v+1 / 2|\mathbf{u}|^{2}\right)(\mathbf{k} \times \mathbf{m})|\mathbf{m}|^{-3}
\end{gathered}
$$

by virtue of which it is possible to determine $u$ in terms of the known quantities $h, a, b, p, v, x_{p}$,
$\mathbf{x}_{\gamma}, t_{\beta}, t_{\gamma}$. Then, $n$ is determined using formula (2.2), and $D_{n}=v_{n}+(\mathbf{u n})$. At points $\gamma_{0}$ the normal to the rigid wall coincides with vector $m|m|^{-1}$, hence $h=0$ on $\gamma_{0}$. It is, consequently, possible to determine on $\gamma_{0}$ all parameters of gas behind the front, the normal $n$, and $D_{n}$ as analytic functions of variables $\beta, \gamma$ at those points $\gamma_{\theta}$, where the regular reflection ( $p<p_{*}$ ) is possible.
3. Transformation of equations and boundary conditions. We introduce in region $\Omega_{3}$ new coordinates, and the functions that define the passage to them are determined in conformity with Sect. 2 as the solution of the Cauchy problem

$$
\begin{aligned}
& y_{\alpha}=D_{n}\left(\left|\mathbf{y}_{\beta} \times \mathfrak{s}_{\gamma}\right|^{2}-\left.D_{n}^{2}\left|T_{\beta} \mathbf{y}_{\gamma}-T_{y_{k}}\right|^{2}\right|^{2}\right)^{-1 / \alpha}\left(\mathbf{y}_{\beta} \times \mathbf{y}_{\gamma}\right), \quad T_{\alpha}=1 \\
& \left.\mathrm{y}\right|_{\alpha=0}=\mathrm{x}_{0}(\beta, \gamma),\left.\quad T\right|_{\alpha=0}=t_{0}(\beta, \gamma) \\
& \mathbf{x}_{\tau}=h\left(\mathbf{x}_{\beta} \times \mathbf{x}_{\mathbf{y}}\right), \quad t_{\tau}=1,\left.\quad \mathbf{x}\right|_{\tau=0}=\left.\mathbf{y}\right|_{\tau=0},\left.\quad t\right|_{\tau=0}=T
\end{aligned}
$$

where $D_{n}$ is a function of $h, \mathbf{y}_{\beta}, \mathbf{y}_{v}, T_{\beta}, T_{v}$ and of quantities ahead of the front considered to be functions of $y, T$, as determined in Sect.2. For the transformation of equations of molionwe use the formulas

$$
\begin{aligned}
& \nabla=\lambda \frac{\partial}{\partial \tau}+\mu \frac{\partial}{\partial \alpha}+I^{-1}\left\{\left(\mathbf{x}_{v}-t_{\gamma} \mathbf{x}_{\alpha}\right) \times\left(\mathbf{x}_{\mathrm{\tau}}-\mathbf{x}_{\alpha}\right) \frac{\partial}{\partial \beta}+\left(\mathbf{x}_{\tau}-\mathbf{x}_{\alpha}\right) \times\left(\mathbf{x}_{\beta}-t_{\beta} \mathbf{x}_{\mathrm{t}}\right) \frac{d}{\partial \gamma}\right\} \\
& \frac{\partial}{\partial t}=J \frac{\partial}{\partial \tau}+|\mathbf{m}|^{-2}\left\{\mathbf{u}\left(\mathbf{x}_{\mathrm{v}} \times \mathbf{m}\right) \frac{\partial}{\partial \beta}-\mathbf{u}\left(\mathbf{x}_{\beta} \times \mathbf{m}\right) \frac{\partial}{\partial \gamma}\right\} \\
& \lambda=I^{-1}\left(\mathbf{x}_{\beta}-t_{\beta} \mathbf{x}_{\alpha}\right) \times\left(\mathbf{x}_{v}-t_{\gamma} \mathbf{x}_{\alpha}\right) . \\
& \mu=-I^{-1}\left(\mathbf{x}_{\beta}-t_{\beta} \mathbf{x}_{\mathbf{x}}\right) \times\left(\mathbf{x}_{\mathrm{y}}-t_{\mathbf{v}_{\mathbf{x}}}\right) . \\
& I=\left(\mathbf{x}_{\mathrm{v}} \quad \mathbf{x}_{\alpha}\right)\left\{\left(\mathbf{x}_{\mathrm{p}}-t_{\beta} \mathbf{x}_{\mathrm{l}}\right) \times\left(\mathbf{x}_{\mathrm{v}}-t_{\gamma} \mathbf{x}_{\mathrm{\tau}}\right)\right\}, \quad J=(\mathbf{q} \mathbf{m})|\mathbf{m}|^{-2}
\end{aligned}
$$

The feasibility of passing to new variables depends on the quantity $I$ not becoming zero or infinity. When $\tau \ldots 0$ we have

$$
I=-|\mathbf{m}|^{-2} v_{n} J^{-1}\left(|\mathbf{m}|^{2}-I_{n}^{2}|\mathbf{k}|^{2}\right)^{-1 / 2}
$$

This formula defines the required property of $I$, at least on $\gamma_{0}$.
The transformed equations (1.2) are schematically represented in the form

$$
\begin{align*}
& s_{\tau}=d_{1} s_{\beta}+d_{2} s_{\gamma} \quad d_{\tau}=e_{1} \mathbf{V}_{\beta}+e_{2} \mathbf{V}_{\gamma}, \quad b_{\tau}=e_{3} \mathbf{V}_{\beta}+e_{\mathrm{s}} \mathbf{V}_{v}  \tag{3.2}\\
& \mathbf{U}_{\boldsymbol{\tau}}=B \mathbf{U}_{\alpha}+D \mathbf{U}_{\alpha}+E_{1} \mathbf{V}_{\beta}+E_{\mathbf{2}} \mathbf{v}_{\gamma}+H_{1}\left(\mathbf{x}_{\alpha}\right)_{\beta}+H_{\mathbf{2}}\left(\mathbf{x}_{\alpha}\right)_{\gamma} \\
& \mathrm{U}=\| \|_{p}^{h} \| \text {, }
\end{align*}
$$

where $\mathbf{V}$ is the vector solution whose components $s, a, b, h, p, x_{\beta}, x_{\gamma}, y_{\beta}, y_{\nu}, z_{\beta}, z_{\gamma}$ and coefficients $d_{i}, e_{i}, B, D, E_{i}, I_{i}$ are scalar, vector-valued, and matrix-valued analytic functions of variables $\mathbf{V}, \mathbf{x}_{\alpha}, \beta, \gamma, \alpha$. Matrix $D$ vanishes when $\tau=0$.

In conformity with sect. 2 boundary conditions are of the form

$$
\begin{align*}
& \tau=0: \quad s=f_{1}\left(h, \mathbf{x}_{\beta}, \mathbf{x}_{\gamma}, \mathbf{x}, \alpha, \beta, \gamma\right), \quad a=f_{2}\left(\mathbf{x}_{\beta}, \mathbf{x}, \alpha, \beta, \gamma\right)  \tag{3.3}\\
& b=f_{3}\left(\mathbf{x}_{\gamma}, \mathbf{x}, \alpha, \beta, \gamma\right), \quad p=f_{1}\left(h, \mathbf{x}_{\beta}, \mathbf{x}_{\gamma}, \mathbf{x}, \alpha, \beta, \gamma\right) ; \quad \alpha=0: \quad h=0
\end{align*}
$$

where $i_{i}$ are analytic functions of their arguments close to the values of the latter on $\gamma_{0}$. The dependence of $f_{i}$ on $x, \alpha, \beta, \gamma$ is linked to the appearance in the relations at the shock wave front of parameters of gas ahead of the front, which are known as functions of $\mathbf{x}, 7$. The input problem is, thus, reduced to solving problem (3.1)-(3.3) followed by the inversion of representations that define transition to new variables.
4. Construction of solution in the class of formal power series. The problem formulated here is a generalization of the Goursat problem. The problem of derivation of its solution in the class of formal power series in variables $\tau, \alpha, \beta-\beta_{1}, \gamma-\gamma_{1}$, where ( $0,0, \beta_{1}, \gamma_{1}$ ) are the coordinates of an arbitrary point $N$ of surface $\gamma_{0}$, reduces to the calculation of derivatives of the solution at that point.

We write the last equation of system (3.2) in the form

$$
\begin{equation*}
\mathbf{C}_{\tau}=B_{N} \mathbf{C}_{\alpha} \mid \mathbf{F} \quad\left(B_{N}=\left.B\right|_{\tau=\alpha=0, \beta, \beta=\beta, \gamma=\gamma, \gamma}\right) \tag{4.1}
\end{equation*}
$$

whose two correlations

$$
\begin{gather*}
B_{N}^{j} \mathbf{U}_{j, n-j}-\mathbf{U}_{0, n}+\sum_{k=1}^{j} D_{N}^{k-1} \mathbf{F}_{k-1, n, n-k}=\sum_{k=1}^{j} B_{N}^{k-1}\left(B_{N} \mathbf{U}_{1,0}-\mathbf{U}_{0,1}+\mathbf{F}_{k-1, n-k}=0\right.  \tag{4.2}\\
B_{N}^{j-n} \mathbf{U}_{j, n-j}-\mathbf{U}_{n, 0}-\sum_{k=j+1}^{n} B_{N}^{k-1-n} \mathbf{F}_{k-1, n-k}=\sum_{k=j+1}^{n} B_{N}^{k-1-n}\left(\mathbf{U}_{0,1}-B_{N} \mathbf{U}_{1,0}-\mathbf{F}\right)_{k-1, n-k}=0 \\
\left.\left(f_{i, j}=\partial^{(i+j)}\right) / \partial \alpha^{i} \partial \tau^{j}\right)
\end{gather*}
$$

are useful for calculating the derivatives.
If $\Lambda_{1} \mathrm{U}_{0, n}$ is specified for $\alpha=0$, and for $\tau=0$ is $\boldsymbol{\Lambda}_{2} \mathrm{U}_{n, 0}$, then by multiplying (4.2) by $\boldsymbol{\Lambda}_{1}, \Lambda_{2}$, we respectively obtain a system of two equations in $\mathbf{U}_{j, n-1}$. In the considered here case $\Lambda_{1}=(1,0), \Lambda_{2}=(-x, 1)\left(x=\partial f_{4} /\left.\partial h\right|_{\tau=0}\right)$. To clarify the question of the obtained system solvability, it is expedient to reduce matrix $B_{N}$ to the diagonal form. The eigenvalues of matrix $B_{N}$ are of the form

$$
v_{1,2}=\left.\left\{(\lambda \mu) \pm-\left(J^{2}|\mu|^{2} c^{-2}-|\lambda \times \mu|^{2}\right)^{1 / 4}\right\}\left\{J^{2} c^{-2}-|\lambda|^{2}\right\}^{-1}\right|_{N}
$$

and real when $\left|w_{N}\right| \geqslant c_{N}$ (subscript $N$ denotes quantities at point $N$ ). Below, the case of strict inequality is called supersonic, and that of strict reverse inequality, subsonic.

In the neighborhood of point $Q$ on $\gamma_{0}$ (i.e. for small $t$ ) we always have the supersonic case, while for large $t$ transition to the subsonic case is possible (but not obligatory) depending on the geometry of incipient wave, its velocity, etc. Since all quantities on roare in conformity with Sect.2, determined prior to the solution of the problem, the "subsonic" and "supersonic" points, as well as the points at which the limit angle of turn of vector $w$, are a priori known.

In the subsonic case $v_{1}=\bar{v}_{2}$. For $v_{1} \neq \nu_{2}$ matrix $B_{N}$ reduces to the diagonal form

$$
B_{N}=M^{-1} \chi M, \quad \chi=\operatorname{diag}\left\{v_{2} ; v_{1}\right\}
$$

It is convenient to introduce new unknown functions

$$
\begin{aligned}
& \mathbf{r}=\|r\|=\mu \mathbf{U}=\left\|\begin{array}{r}
h+K p \\
-h+K p
\end{array}\right\| \\
& K=-I\left(J^{2}|\boldsymbol{\mu}|^{2} c^{-2}-|\lambda \times \boldsymbol{\mu}|^{2}\right)^{1 /} \rho^{-1}(J|\mathbf{m}|)_{N}^{-2}
\end{aligned}
$$

Equations (4.1) assume the form

$$
\begin{equation*}
r_{\tau}=v_{2} r_{\alpha}-\psi^{1}, \quad l_{\tau}=v_{1} l_{\alpha}+\psi^{2}, \quad\left\|-\psi^{2}\right\|=. / / F \tag{4.3}
\end{equation*}
$$

Solving the equations obtained by multiplying (4.2) by $\boldsymbol{\Lambda}_{1}, \boldsymbol{\Lambda}_{2 p}$ we obtain the formulas

$$
\begin{align*}
& r_{j, n-j}=\Delta_{n}^{-1}\left\{\sum_{k=1}^{n} v_{2}^{n-j} v_{1}^{k-1-n} \psi_{k-1, n-k}^{2}+\sum_{k=1}^{j} v_{2}^{n-j+k-1} v_{1}^{-n} \psi_{k-1, n-k}^{1}+\right.  \tag{4.4}\\
& \left.\sum_{k=j+1}^{n} d v_{2}^{k-j-1} \psi_{k-1, n-k}^{1}+v_{2}^{n-j} v_{1}^{-n}\left(r_{0, n}-l_{0, n}\right)+v_{2}^{n-j}\left(l_{n, 0}-d r_{n, 0}\right)\right\} \\
& l_{j, n-j}=\Delta_{n}^{-1}\left\{\sum_{k=1}^{n} d v_{2}^{k-1} v_{1}^{-j} \psi_{k-1, n-k}^{\prime}+\sum_{k=1}^{j} d v_{1}^{k-j-1} \psi_{k-1, n-k}^{2}+\right. \\
& \left.\sum_{k=j+-1}^{n} v_{2}^{n} v_{1}^{k-n-j-1-1} \psi_{k-1, n-k}^{3}+d v_{1}^{-j}\left(r_{0, n}-l_{0, n}\right)+v_{2} v_{1}^{-j}\left(l_{n, 0}-d r_{n, 0}\right)\right\} \\
& \Delta_{n}=v_{2}^{n} v_{1}^{-n} \quad d, \quad d=(\kappa K-1)(\kappa K \quad 1)^{-1}
\end{align*}
$$

These formulas are used for calculating the derivatives of order $n$ with respect to variables $\tau, \alpha$ of functions $r, l$ in terms of derivatives of order $n-1$ of $\psi^{i}$ and derivatives $\left(r_{0, n}-l_{0, n}\right)_{\alpha=0},\left(l_{n, 0}-d r_{n, 0}\right)_{\tau=0}$ are known by virtue of (3.3). The necessary condition of solvability of equations for $r_{j}, n-j, l_{j, n-j}$ are of the form $\Delta_{n} \neq 0, n=1,2, \ldots$ If $\Delta_{n}$ vanishes, we obtain a contradiction to the assumption of existence an analytic solution of the problem with arbitrary data.

In the supersonic case $v_{2}>v_{1}>0,|d|<1$ and the conditions of solvability are always satisfied. In the subsonic case $v_{2}=\bar{v}_{1}$, and $K$ is an imaginary number. If $v_{2} v_{1}^{-1}=c^{2 \pi i \sigma}$, $d=e^{2 \pi i \delta} \quad\left(0 \leqslant \sigma \leqslant 1,0 \leqslant \delta \leqslant 1\right.$, the solvability condition is of the form $e^{2 \pi i n \sigma} \ldots c^{2 \pi i \delta} \neq 0$ for all positive integers $n$. The quantities $J_{n}$ vanish on straight lines $\sigma=\delta n^{-1}+j^{-1}(j=0,1$, . .., $n-1$ ) that lie in the unit square of the plane ( $\sigma, \delta$ ). Points of these straight lines for various $n, j$ form everywhere in that square a compact set. Then in any neighborhood of point $(\sigma, \delta)$ with $\Delta_{n} \neq 0(n=1,2, \ldots)$ there are points at which certain $\Lambda_{k}=0$. This property can be interpreted as instability of the subsonic problem in the class of analytic functions. Since on $\gamma_{0}, \sigma$ and $\delta$ are continuous functions of $\beta$ and $\gamma$, it is possible to maintain that when $\sigma \not \equiv$ const, $\delta \neq$ const there are points on $\gamma_{0}$ at which $\Delta_{h}=0$. Consequently the problem of regular reflection has, generally, no analytic solution in the subsonic case.

Below, we consider the supersonic case. Jt is convenient to transform the boundary conditions that link $r$ and $l$ to the form

$$
\begin{equation*}
(r-l)_{\alpha=0}=0, \quad(l-d r)_{\tau=0}=g_{1}\left(r, \mathbf{x}_{\beta}, \mathbf{x}_{\gamma}, \mathbf{x}, \alpha, \beta, \gamma\right) \tag{4.5}
\end{equation*}
$$

where $g_{1}$ is an analytic function of its arguments, such that $\left(\sigma g_{1} / \partial r\right)_{N}=0$.

Lemma. A unique solution of problem(3.1)-(3.3) exists in the class of formal power series in variables $\tau, \alpha, \beta-\beta_{v}, \gamma-\gamma_{1}$.

Proof is by induction over the total number $n$ of differentiation with respect to $\tau$ and $\alpha$. In accordance with sect. 2 all sought functions are known on $\gamma$ when $\tau=\alpha=0$. Differentiating them with respect to $\beta$ and $\gamma$, we obtain all derivatives with respect to these variables. If all derivatives with the over-all $n-1$ of differentiation withrespect to $\tau, a$, the derivatives of $r$ and $l$ of order $n$ with respect to $\tau$ and $\alpha$ are determined using (4.4) and the formulas obtained by differentiating (4.4) with respect to $\beta$ and $\gamma$. Derivatives of remaining functions are determined by equations obtained by differentiating (3.1)-(3.3).
5. Existence of analytic solution of the problem in the supersonic case. Convergence of formal power series is proved by costructing majorants of solutions that are obtained by solving the auxiliary majorant problem. As a preliminary, the nonlinear equations (3.1) in $\mathbf{x}$ and $\mathbf{y}$ are reduced to quasi-linear equations in $\mathbf{x}$ and $\mathbf{y}$ by extending them to derivatives, and boundary conditions of the problem are reduced to a homogeneous form by substituting the unknown functions (subscript unity denotes new notation) $y_{1}=\mathrm{y}-\mathrm{x}_{0}(\beta, \gamma), \mathrm{x}_{1}=\mathrm{x}-\mathrm{y}$ and the respective substitutions of derivatives of these functions; further

$$
s_{1}=s-f_{1}, \quad a_{1}=a-f_{2}, \quad b_{1}=b-f_{3}, \quad r_{1}=i-g_{1}(1-d)^{-1}, \quad l_{1}=l-g_{1}(1-d)^{-1}
$$

For the transformation of quantities we obtain the homogeneous boundary conditions

$$
\begin{aligned}
& z=0: \quad \mathbf{y}_{1}=\mathbf{y}_{1 \beta}=\mathbf{y}_{1 \gamma}=0, \quad r_{1}-l_{1}=0 \\
& \tau=0: \quad \mathbf{x}_{1}=\mathbf{x}_{1 \beta}=\mathbf{x}_{1 \gamma}=\mathbf{x}_{1 \alpha}=0, \quad a_{1}=b_{1}=s_{1}=l_{1}-d r_{1}=0
\end{aligned}
$$

Formulas of the form (4.4) also apply to transformed $\psi^{1}$ and $\psi^{2}$. These formulas and the transformed equations imply that the problem with boundary conditions of form (5.1), where the equality sign is replaced by the majorizing relation $\left(\left.s_{1}\right|_{\tau=0} \gg 0\right.$, etc. $)$ can be taken as the majorizing problem. We recall that the relation $f>g$ means that the coefficients of expansion of function $f$ in series in powers of its arguments are not lower than the absolute values of respective coefficients of $g$.

Equations of the majorizing problem are obtained as follows. Coefficients at derivatives in the right-hand sides of transformed equations and in the expressions of transformed functions $\psi^{1}, \psi^{2}$ are replaced by their majorants, with the retention of the property of some coefficients vanishing at point $N$. Functions that satisfy "like" equations are majorized by one majorant.

Let $Y$ be the majorant for all components $y_{1}, y_{1 \beta}, y_{1 \gamma}, X$ the majorant for components $x_{1}$, $x_{1 \beta}, x_{1 v}, Z$ the majorant for components $x_{1 \alpha}, S$ the majorant for function $s_{1}, R$ and $L$ the majorants for $r_{1}$ and $l_{1}$, and $A$ the majorant for $a_{1}$ and $b_{1}$. The majorant system is of the form

$$
\begin{align*}
& Y_{\alpha}=F_{1} \Psi_{1}, \quad X_{\tau}=F_{2}\left(Y_{\tau} \ldots \Psi_{1}\right), \quad Z_{\tau}=F_{3}\left(L_{\alpha}+R_{\tau}+\Psi_{1}+\Psi_{2}\right)  \tag{5.2}\\
& S_{\tau}=F_{4}\left(L_{\alpha}-R_{\tau}+\Psi_{1}+\Psi_{2}\right), \quad \Lambda_{\tau}=F_{\mathbf{5}} \Psi_{1}, \\
& \gamma_{2} R_{\alpha}=R_{\tau}+F_{6}\left(\Psi_{1}+\Psi_{2}\right) \\
& L_{\tau}=v_{1} L_{\alpha}+F_{6}\left(\Psi_{1}+\Psi_{2}\right), \quad \xi=\eta_{1} \tau-\eta_{2} \alpha+\beta-\beta_{1}+\gamma-\gamma_{1}, \\
& \Psi_{1}=\Sigma_{\beta}+\Sigma_{\gamma} \div 1, \quad \Psi_{2}=(\Sigma+\xi)\left(\Sigma_{\tau}+\Sigma_{\alpha}\right), \Sigma=(Y+X+Z+S+A+R+L)
\end{align*}
$$

where $\eta_{1}$ and $\eta_{2}$ are constants which will be defined later, and $F_{i}$ are the majorants of coefficients of input equations.

The solution of system (5.2) must vanish at point $N$. Hence the majorants of coefficients in that point neighborhood can be taken in the form

$$
F_{i}=K_{i}\left(1-(\Sigma-\xi) M_{i}\right)^{-1}
$$

with appropriately selected constants $K_{i}, M_{i}$. Formulas similar to (4.4) can be obtained for Eqs.(5.2), from which follows that, if we set in boundary conditions (5.1) $d=1$, the obtained problem is also a majorant for the input problem.

We seek a particular solution of Eqs.(b.2) that depends only on $\xi$, vanishes at point $\xi=0$, and such that $R \equiv L$. By virtue of (5.2) the last requirement provides the following relation between $\eta_{1}$ and $\eta_{2}: \eta_{1}=1_{2}\left(v_{2}+v_{1}\right) \eta_{2}$. The system of equations for finding the particular solution $X(\xi), Z(\xi), S(\xi), A(\xi), Y(\xi), H(\xi) \equiv L(\xi)$ is of the form

$$
\begin{aligned}
& \eta_{2} Y^{\prime}=F_{1}\left(2 \Sigma^{\prime}+1\right), \quad x_{1} \eta_{0} X^{\prime}=F_{2}\left(\varkappa_{1} n_{2} Y^{\prime}+2 \Sigma^{\prime}+1\right) \\
& x_{1} \eta_{2} Z^{\prime}=F_{3}\left(\eta_{2} L^{\prime}+x_{1} \eta_{2} R^{\prime}+\Psi_{3}\right), \quad x_{1} \eta_{2} A^{\prime}=F_{5}\left(2 \Sigma^{\prime}+1\right) \\
& x_{1} \eta_{2} S^{\prime}=F_{4}\left(\eta_{2} L^{\prime}+x_{1} \eta_{2} R^{\prime}+\Psi_{3}\right), \quad x_{2} \eta_{2} R^{\prime}=F_{8} \Psi_{3} \\
& x_{1}=1 l_{2}\left(v_{2}-v_{1}\right), x_{2}=1_{2}\left(v_{2}-v_{3}\right), \quad \Psi_{3}=\left((\Sigma+\xi) \eta_{2}\left(1+x_{1}\right)+2\right) \Sigma^{\prime}+1
\end{aligned}
$$

The substitution of expressions for $Y^{\prime}, R^{\prime} . L^{\prime}$ from the first and last of equations into the right-hand sides of the remaining ones together with the linear combination of equations yield for function $\Sigma$ the equation

$$
\begin{aligned}
& \eta_{2} \Sigma^{\prime}=\left(F_{1}+x_{1}^{-1} F_{2}\left(x_{1} F_{1}+1\right)+x_{1}^{-1} F_{\mathrm{b}}\right)\left(2 \Sigma^{\prime}+1\right)+ \\
& \quad\left(x_{1}^{-1}\left(F_{3}+F_{4}\right)\left(\left(1+x_{1}\right) x_{2}^{-1} F_{6}-1\right) \div 2 x_{2}^{-1} F_{6}\right)\left(\Psi_{3}-1\right)
\end{aligned}
$$

In the neighborhood of point $\Sigma=\xi=0$ this equation can be solved for $\Sigma^{\prime}$. For this a fairly large parameter $\eta_{2}$ is selected, namely

$$
\begin{aligned}
& \eta_{2}>2\left(K_{1}+x_{1}^{-1} K_{2}\left(x_{1} K_{1}+1\right)+x_{1}^{-1} K_{5}\right)+ \\
& \quad 2\left(x_{1}^{-1}\left(K_{3}+K_{4}\right)\left(\left(1 \because x_{1}\right) x_{2}^{-1} K_{6}+1\right)+2 x_{2}^{-1} K_{6}\right)
\end{aligned}
$$

The equation for the determination of $\Sigma$ now assumes the form

$$
\begin{equation*}
\Sigma^{\prime}=\Phi(\Sigma, \xi) \tag{5.3}
\end{equation*}
$$

with function $\Phi$ of the majorant type (with positive coefficients of expansion in series in powers of its arguments in the neighborhood of point $\Sigma=\xi=0$ ). By the Cauchy - Kovalevsky theorem there exists an analytic solution of Eq. (5.3) that is, also, a function of the majorant type, and $\Sigma(0)=0$. Using the known function $\Sigma$ we determine functions $Y, X, Z, S, A, R=L$ which are also of the majorant type. We, thus, obtain the particular solution of the majorant systom (5.2), which satisfies the conditions

$$
\left.Y\right|_{\alpha=0} \gg 0,\left.\quad X\right|_{\tau=0} \gg 0,\left.\quad Z\right|_{\tau=0} \gg 0,\left.\quad S\right|_{r=0} \gg 0,\left.\quad A\right|_{\tau=0} \gg 0, \quad R=L
$$

A successive calculations of coefficients of expansion of this sulution in power series will show that these series majorize the power series constructed in Sect.4.
6. The one-sheeted mapping of $(\tau, \alpha, \beta, \gamma) \rightarrow(t, x, y, z)$. The Jacobian of that mapping is finite and nonzero at points of $\gamma_{0}$. This enables us to use locally in the neighborhood of every point of $\gamma_{0}$ the implicit function theorem. For small $\tau, \alpha$ (i.e. in the $\gamma_{0}$ neighborhood) the mapping is one-sheeted, hence the equalities

$$
t\left(\tau_{1}, \alpha_{1}, \beta_{1}, \gamma_{1}\right)=t\left(\tau_{2}, \alpha_{2}, \beta_{2}, \gamma_{2}\right), \quad \mathbf{x}\left(\tau_{1}, \alpha_{1}, \beta_{1}, \gamma_{1}\right)=\mathbf{x}\left(\tau_{2}, \alpha_{2}, \beta_{2}, \gamma_{2}\right)
$$

imply the equalities

$$
\tau_{1}=\tau_{2}, \alpha_{1}-\alpha_{2}, \quad \beta_{1}-\beta_{2}, \quad \gamma_{1}-\gamma_{2}
$$

By the definition of function $t$ we have $\alpha_{1}+\tau_{1}+t_{0}\left(\beta_{1}, \gamma_{1}\right)=\alpha_{2}+\tau_{2}+t_{0}\left(\beta_{2}, \gamma_{2}\right)$. In conformity with (3.1) we have for functions $\mathbf{x}(\tau, \alpha, \beta, \gamma)$ valid the following formula

$$
\begin{equation*}
\mathbf{x}(\tau, \alpha, \beta, \gamma)=x_{0}(\beta, \gamma)+\int_{0}^{\alpha} \mathbf{H}(0, \xi, \beta, \gamma) d \xi+\int_{0}^{\tau} \mathbf{G}(\eta, \alpha, \beta, \gamma) d \eta \tag{6.1}
\end{equation*}
$$

where $\mathbf{H}$ and $\mathbf{G}$ coincide with the right-hand sides of Eqs.(3.1), and it is possible to substitute in $H$ function $x$ for $y$ owing to their equality when $\tau=0$.

Consider the neighborhood of $\gamma_{0}$ in which $|H|<C_{1},|G|<C_{1}\left(C_{1}\right.$ is a positive constant $)$. From (6.1) we have the inequality

$$
\left|\mathbf{x}(\tau, \alpha, \beta, \gamma)-\mathbf{x}_{\mathrm{n}}(\beta, \gamma)\right|<c_{1}(\tau-\alpha \alpha)
$$

Then the following two inequalities are also satisfied:
$\left|\mathbf{x}_{0}\left(\beta_{1}, \gamma_{1}\right)-\mathrm{x}_{0}\left(\beta_{2}, \gamma_{2}\right)\right|<C_{1}\left(\tau_{1}+\alpha_{1}+\tau_{2}+\alpha_{2}\right\rangle,\left|t_{0}\left(\beta_{1}, \gamma_{1}\right)-t_{0}\left(\beta_{2}, \gamma_{2}\right)\right|<\tau_{1}+\alpha_{1} \div \tau_{2}+\alpha_{2}$
The specified mapping $t=t_{0}(\beta, \gamma), \mathbf{x}=\mathbf{x}_{0}(\beta, \gamma)$ is such that from the last inequalities we have

$$
\left|\beta_{1}-\beta_{2}\right|+\left|\gamma_{1}-\gamma_{2}\right|<c_{2}\left(\tau_{1}+\alpha_{1}+\tau_{2}+\alpha_{2}\right)
$$

with some positive constant $C_{2}$. By virtue of (6.1) the relation $x\left(\tau_{1}, \alpha_{1}, \beta_{1}, \gamma_{1}\right)=x\left(\tau_{2}, \alpha_{2}, \beta_{2}, \gamma_{2}\right)$ can be written in the form

$$
\begin{aligned}
& \mathbf{x}_{0}\left(\beta_{1}, \gamma_{1}\right)-\mathbf{x}_{0}\left(\beta_{2}, \gamma_{2}\right)-\left(\alpha_{2}-\alpha_{1}\right) \mathbf{H}\left(0,0, \beta_{2}, \gamma_{2}\right)- \\
& \quad\left(\tau_{2}-\tau_{1}\right) \mathbf{G}\left(0, \alpha_{2}, \beta_{2}, \gamma_{2}\right)=-\int_{0}^{\alpha_{1}}\left(\mathbf{H}\left(0, \xi, \beta_{1}, \gamma_{1}\right)-\mathbf{H}\left(0, \xi, \beta_{2}, \gamma_{2}\right)\right) d \xi- \\
& \int_{0}^{\tau_{1}}\left(\mathbf{G}\left(\eta, \alpha_{1}, \beta_{1}, \gamma_{1}\right)-G\left(\eta, \alpha_{2}, \beta_{2}, \gamma_{2}\right)\right) d \eta+\int_{\alpha_{1}}^{\alpha_{2}}\left(\mathbf{H}\left(0, \xi_{\xi}, \beta_{2}, \gamma_{2}\right)-\right. \\
& \left.\mathbf{H}\left(0,0, \beta_{2}, \gamma_{2}\right)\right) d \xi+\int_{\tau_{1}}^{\tau_{2}}\left(G\left(\eta, \alpha_{2}, \beta_{2} \gamma_{2}\right)-G\left(0, \alpha_{2}, \beta_{2}, \gamma_{2}\right)\right) d \eta
\end{aligned}
$$

from which we obtain the inequality

```
\(\mid\left(x_{c \beta}\left(\beta_{2}, \gamma_{2}\right)-t_{\beta}\left(\beta_{2}, \gamma_{2}\right) G\left(0,0, \beta_{2}, \gamma_{2}\right)\right)\left(\beta_{1}-\beta_{2}\right)+\left(x_{0 \gamma}\left(\beta_{2}, \gamma_{2}\right)-\right.\)
    \(\left.\boldsymbol{t}_{\boldsymbol{y}}\left(\beta_{2}, \gamma_{2}\right) \mathbf{G}\left(0,0, \beta_{2}, \gamma_{2}\right)\right)\left(\boldsymbol{\gamma}_{1}-\gamma_{2}\right) \cdots\left(\mathbf{H}\left(0,0, \beta_{2}, \gamma_{2}\right)-\right.\)
    \(\left.\mathbf{G}\left(0,0, \beta_{2}, \gamma_{2}\right)\right)\left(\alpha_{1}-\alpha_{2}\right) \mid \leqslant C_{3}\left(\tau_{1} \div \alpha_{1} \div \tau_{2}+\alpha_{2}\right)\left(\mid \beta_{1}-\right.\)
    \(\beta_{2}\left|\div\left|\gamma_{1}-\gamma_{1}\right|+\left|\alpha_{1}-\alpha_{2}\right|\right)\)
```

with some constant $c_{3}>0$. From the obtained inequality we obtain for small $\tau_{1}, \tau_{2}, \alpha_{1}, \alpha_{8}$ the required equalities by virtue of the linear independence of vectors $\mathbf{x}_{\beta}-t_{\beta} \mathbf{G}, \mathbf{x}_{\gamma}-t_{\gamma} \mathbf{G}, \mathbf{H}-\mathbf{G}$ at points of $\gamma_{0}$. The one-sheeted property of mapping is proved. In the indicated neighborhood of $\gamma_{n}$ the mapping can be inverted yielding a solution in the form of analytic functions of $x, t$.

The region in which the problem was solved compriscs the band $0 \ll t_{0}$ with some $t_{1}>0$. The obtained solution, thus, defines the initial stage of shock wave reflection from the wall. An analytic solution of the problem of regular reflection that is uniquely determined in some neighoorhood of the incident wave trace moving along the wall was obtained for large i under the condition of analyticity of function that define the gas flow in the neighborhood of the incident wave front. The above reasoning shows that, when the projection on the $\Pi$ plane of gas velocity relative to the trace behind the reflected shock wave front becomes subsonic (at small $t$ it is higher than the local speed of sound), the conditions of solvability of the problem in the class of analytic functions are not satisfied, which implies the appearance of solution singularities. In the case of convergence of the derived series and one-to-one mapping $(\tau, \alpha, \beta, \lambda) \rightarrow(t, x, y, z)$ the obtained solution defines the total flow of gas behind the reflected wave, and not only in the trace neighborhood. The transition from a regular shock wave reflection to the irregular one requires investigation.

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